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ALGEBRAS GROUPS AND GEOMETRIES

VOLUME 28, NUMBER 3, SEPTEMBER 2011

ON COMMUTATIVE WEAKLY SCI – RINGS, 233

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FRACTIONS CONTINUES DE MATRICES
ET DÉVELOPPEMENTS

DE CERTAINES FONCTIONS SPÉCIALES, 245

Ali Kacha, Brahim Ounir

Équipe de Géométrie Complexe et

Approximations Diophantiennes

Laboratoire Amsi

Université IBN Tofail

B.P. 133, Kénitra, 14 000, Maroc

QUASI δ -DISCRETE MODULES, 259

Mohammad Javad Nematollahi

Department of Mathematics, Islamic Azad University

Arsanjan Branch, P.O. Box 73761-168, Arsanjan, Iran

TOPOLOGICAL COLIMIT PRESERVING EQUALIZERS, 275

N. Hosseini

Dept. of Math., Shahid Bahonar University of Kerman, Iran

Gh. Mirhosseinkhani

Dept. of Math., Hormozgan University, Bandarabas, Iran

**UNIFICATION OF q -FRACTIONAL
INTEGRAL OPERATORS, 283**

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Facultad de Ingeniería. Universidad del Zulia
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**STUDIES ON THE LIE-SANTILLI ISOTHEORY
WITH UNIT OF GENERAL FORM, 299**

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ON COMMUTATIVE WEAKLY SCI - RINGS**Mamadou Barry and Papa Cheikhou Diop**

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Abstract

Let R be a commutative ring and M an unital R -module. M is called co-Hopfian if any injective endomorphism of M is an isomorphism. M is called weakly co-Hopfian if any injective endomorphism of M is essential. The ring R is called SCI-ring if any co-Hopfian R -module is finitely cogenerated. R is called weakly SCI-ring if any weakly co-Hopfian R -module is finitely cogenerated. In this note we show that a commutative ring R is a weakly SCI-ring if and only if it is an Artinian principal ideal ring.

Keywords : SCI-ring, weakly SCI-ring, Artinian principal ideal ring, finitely generated module, finitely cogenerated module, co-Hopfian module, weakly co-Hopfian module, uniform dimension.

Classification: Mathematical subject classifications: 13E10, 13C05.

Throughout all rings are associative commutative with identity and all modules are unital module. Let R be a ring and let M be an R -module. A submodule L of M is called essential if $L \cap N \neq 0$ for every non-zero submodule N of M . The socle of M , denoted by $Soc(M)$, is the sum of the zero submodule and all simple submodules of M . M is said to have finite uniform dimension if it does not contain a direct sum of an infinite number of non-zero submodules. It is known that for an R -module M with a finite uniform dimension, there exists a positive integer n such that every direct sum of non-zero submodule of M does not have more than n summands. We denote the uniform dimension of M by $udim(M)$. M is called co-Hopfian if any injective endomorphism of M is an isomorphism. M is called weakly co-Hopfian if any injective endomorphism of M is essential. It is well known that any co-Hopfian R -module is weakly co-Hopfian. M is finitely cogenerated if $Soc(M)$ is finitely generated and essential in M . The commutative ring R is called SCI-ring if any co-Hopfian R -module is finitely cogenerated. R is called weakly SCI-ring if any weakly co-Hopfian R -module is finitely cogenerated.

In this paper, we show that for a commutative ring R , the following conditions are equivalent :

- 1 - R is a weakly SCI-ring.
- 2 - R is an Artinian principal ideal ring.

Definitions and notations used in this paper can be found in [1], [10] and [15].

1 Preliminaries

Lemma 1.1. *Let R be a ring and let M be a left R -module. The following assertions are verified :*

1. *If M is co-Hopfian, then M is weakly co-Hopfian.*
2. *If M has finite uniform dimension, then M is weakly co-Hopfian.*
3. *If M is finitely cogenerated, then M has finite uniform dimension.*
4. *If M is finitely cogenerated, then M is weakly co-Hopfian.*

Proof

1. Is obvious.
2. It results from ([7] , proposition, 4.12)
3. It results from ([10], page 537)
4. It results from 2. and 3.

Lemma 1.2. ([10], page 537)

Let R be a ring and let M be a left R -module. If M is finitely cogenerated, then every submodule of M is also finitely cogenerated.

Lemma 1.3. ([10], corollary 6.10)

Let R be a ring and let M be a left R -module. If M has finite uniform dimension, then every submodule N of M has also finite uniform dimension and $\text{udim}(N) \leq \text{udim}(M)$. If N is an essential submodule of M , then $\text{udim}(M) = \text{udim}(N)$.

Proposition 1.4. [14]

Let R be a commutative ring. Then every finitely generated R -module is co-Hopfian if and only if every prime ideal of R is maximal.

Proposition 1.5. ([1], proposition 10.18)

Let R be a unitary ring. Then the following conditions are equivalent :

1. R is left Artinian.
2. Every finitely generated left R -module is Artinian.
3. Every finitely generated left R -module is finitely cogenerated.

2 Example of a non finitely cogenerated module which is weakly co-Hopfian over a local Artinian ring whose maximal ideal is not principal

Let R be a commutative local Artinian ring which is a non principal ideal ring. We may suppose without loss of generality that the ring R is local Artinian with Jacobson radical $J(R) = aR + bR$ where $a \neq 0$, $b \neq 0$ and $a^2 = b^2 = ab = 0$. Then by [6], we may write $R = C \oplus bC$ where C is an Artinian local subring of R , with Jacobson radical $J(C) = aC$ where $a \neq 0$ and $a^2 = 0$. Let M be the total ring of fractions of the polynomial ring $C[X]$, σ be the C -endomorphism of the C -module M defined by $\sigma(m) = aXm$ for $m \in M$; $\varphi : R \rightarrow \text{End}_C M$ be the homomorphism of rings defined by $\varphi(\alpha + \beta b) = \alpha 1_M + \beta \sigma$ for $\alpha + \beta b \in R$ where $\alpha \in C$, $\beta \in C$ and 1_M be the identity homomorphism of M . We consider on M the R -module structure defined by $(\alpha + \beta b)m = \varphi(\alpha + \beta b)(m) = \alpha m + \beta aXm$, for $\alpha + \beta b \in R$, ($\alpha \in C$, $\beta \in C$) and for $m \in M$. If f is an R -endomorphism of the R -module M , then for all $m \in M$ we have :

$$\sigma.f(m) = bf(m) = f(bm) = f(\sigma(m)).$$

Thus, the R -endomorphisms of the R -module M are the C -endomorphisms of M commuting with σ . Following ([9] , proposition 2.3) the R -module M is co-Hopfian thus weakly co-Hopfian.

For $d = \lambda a + \gamma b \in J(R) = Ra + Rb$ and for $am \in C[X]$ we have :

$$d.(am) = (\lambda a + \gamma b)am = \varphi(\lambda a + \gamma b)(am) = (\lambda a 1_M + \gamma \sigma)(am) = \lambda a^2 m + \gamma a^2 X m = 0.$$

Therefore, the submodule $aC[X] = \bigoplus_{n \geq 0} aC X^n$ of M is annihilated by $J(R)$, then $aC[X]$ is semi-simple as an R -module. Since $aC[X]$ is not finite length as an R -module, then M is not finitely cogenerated.

3 Characterization of commutative weakly SCI-rings

Proposition 3.1. *Let R be a commutative weakly SCI-ring. If R is an integral domain, then R is a field.*

Proof

Let K be the classical quotient field of the integral domain R . The R -module ${}_RK$ is weakly co-Hopfian. It follows that ${}_RK$ is finitely cogenerated. Thus,

$\text{Soc}({}_RK) \cap_R R \neq 0$. Let $S = Ra$, $a \in R \setminus \{0\}$ a simple submodule of

$\text{Soc}({}_RK) \cap_R R$. The map $\varphi : {}_R R \rightarrow S = Ra$ is an isomorphism of R -modules.

Therefore ${}_R R$ is simple. So for $b \in R \setminus \{0\}$, we have $R = Rb = Rb^2$ then $b = cb^2$ for $c \in R$. It follows that $1 = cb$.

Remark 3.2. :

1. *Any homomorphism image of a commutative weakly SCI-ring is a weakly SCI-ring.*
2. *Direct Product of rings $R_i (1 \leq i \leq n)$ is weakly SCI-ring if and only if so is each R_i for $(1 \leq i \leq n)$.*

Lemma 3.3. *Let an R -module $M = M_1 \oplus M_2$ be a direct sum of submodules M_1 and M_2 . Then M_1 is a fully invariant submodule of M if and only if $\text{Hom}(M_1; M_2) = 0$.*

Proof Suppose first that M_1 is a fully invariant submodule of M .

Let $f : M_1 \longrightarrow M_2$ be any homomorphism.

Let $p_1 : M \longrightarrow M_1$ denote the canonical projection and $i_2 : M_2 \longrightarrow M$ denote the inclusion. Then $f^* = i_2 \circ f \circ p_1$ is an endomorphism of M . By hypothesis, $(i_2 \circ f \circ p_1)(M_1) = f^*(M_1) \subseteq M_1$, so that $f(M_1) \subseteq M_1 \cap M_2 = 0$. It follows that $f = 0$.

Conversely suppose that $\text{Hom}_R(M_1, M_2) = 0$. Let $p_2 : M \longrightarrow M_2$ denote the

canonical projection and $i_1 : M_1 \longrightarrow M$ denote the inclusion. For any endomorphism g of M ,

$$g(M_1) \subseteq p_1 \circ g(M_1) + p_2 \circ g \circ i_1(M_1) = p_1 \circ g(M_1) \subseteq M_1,$$

because

$$p_2 \circ g \circ i_1 \in \text{Hom}_R(M_1, M_2) = 0.$$

It follows that M_1 is a fully invariant submodule of M .

Lemma 3.4. *Let P_1 and P_2 be two primes ideal of R . If $P_1 \not\subseteq P_2$, then $\text{Hom}_R(R/P_1; R/P_2) = 0$.*

Proof Let $f \in \text{Hom}_R(R/P_1, R/P_2)$. We have $f(1 + P_1) = t + P_2$, or $f(\bar{1}) = \bar{t}$, for some $t \in R$. For all $x \in R$, we have $f(\bar{x}) = xf(\bar{1}) = x\bar{t} = \overline{xt}$, or $f(x + P_1) = xt + P_2$. Since f is an R -homomorphism, we have $f(\bar{0}) = \bar{0}$, i.e $f(P_1) \subseteq P_2$. So, let $x \in P_1$ and $x \notin P_2$. We have $f(\bar{x}) = \bar{0}$, thus $xt + P_2 = P_2$. Hence $xt \in P_2$ and $x \notin P_2$ imply that $t \in P_2$, or $f(x + P_1) = 0$.

Lemma 3.5. ([8], proposition 1.11)

Let a module $M = \bigoplus_{i \in I} M_i$ such that each M_i is invariant under any injective endomorphism of M . Then M is weakly co-Hopfian if so is each M_i for $i \in I$.

Proposition 3.6. *Let R be a commutative weakly SCI-ring. Then every prime ideal of R is maximal. Moreover the set of all prime ideals of R is finite.*

Proof

Let P a prime ideal of R , then R/P is a commutative integral domain. Since R/P is homomorphism image of a commutative weakly SCI - ring R , then R/P is a weakly SCI - ring. It follows that R/P is a field and P is maximal.

Let now L be the set of all prime ideals of R and $M = \bigoplus_{P \in L} R/P$.

For any $P \in L$, R/P is a simple R -module, furthermore if P and $P' \in L$ with $P \not\subseteq P'$ then by lemma 3.4, $\text{Hom}_R(R/P; R/P') = 0$, hence R/P is a fully invariant

submodule of M . It follows that for any $f \in \text{End}(M)$, $f = \bigoplus_{P \in L} f_P$ where $f_P \in \text{End}(R/P)$. Since R/P is a simple R -module then f_P is an isomorphism so $f = \bigoplus_{P \in L} f_P$ is an injective endomorphism of M . Then by lemma 3.5 the R -module $M = \bigoplus_{P \in L} A/P$ is weakly co-Hopfian, so $\text{Soc}(M) = M$ is a finitely cogenerated R -module. It follows that M is finitely generated and then the set L is finite.

Corollary 3.7. *Let R be a commutative weakly SCI - ring*

1. *The Jacobson radical $J(R)$ of R is a nil ideal*
2. *The idempotents can be lifted modulo $J(R)$, and R is semi-local.*
3. *R is a finite direct product of commutative local weakly SCI-rings.*
4. *If M is a simple R -module then $M \cong I$ where I is a minimal ideal of R .*

Proof

It follows from proposition 3.6.

Proposition 3.8. *Let R be a commutative weakly - SCI - ring. Then R is Artinian.*

Proof

By proposition 1.5 , it suffices to show that every finitely generated module is finitely cogenerated. Since R is a commutative ring and every prime ideal of R is maximal, then by [14] every finitely generated module is co-Hopfian thus weakly co-Hopfian. From the Hypothesis that R is a weakly SCI-ring it follows that M is finitely cogenerated.

Proposition 3.9. *Let R be a commutative weakly SCI-ring. Then R is a finite direct product of local Artinian weakly SCI-ring.*

Proof It follows from corollary 3.7 and proposition 3.8.

Theorem 3.10. *Let R be a commutative ring. The following conditions are equivalent :*

1. R is a weakly SCI-ring.
2. R is an Artinian principal ideal ring.

Proof

1. $\implies$ 2. If R is a weakly SCI - ring then by proposition 3.9 and §2§ R is an Artinian principal ideal ring.

2. $\implies$ 1. Let R be an Artinian principal ideal. Let M be a weakly co-Hopfian R - module then by ([4], proposition 3.8) M is finitely generated. Since R is Artinian then by ([1], proposition 10.18) M is finitely cogenerated. Therefore R is a weakly SCI - ring.

Corollary 3.11. *Let R be a commutative weakly SCI - ring. The following conditions are equivalent for an R - module M :*

1. M has finite uniform dimension.
2. M is weakly co-Hopfian.
3. $\text{Soc}(M)$ is weakly co-Hopfian.
4. $E(M)$ is weakly co-Hopfian.

Proof

1. $\implies$ 2. is obvious

2. $\implies$ 1. Let M be a weakly co-Hopfian R - module. Since R is a weakly SCI - ring, then M is finitely cogenerated. It follows that M has finite uniform dimension.

1. $\implies$ 3. is obvious

3. $\implies$ 1. Let $\text{Soc}(M)$ be weakly co-Hopfian. Since there are only finitely many simple R - module up to isomorphism then by ([8], corollary 1.12) $\text{Soc}(M)$ has

finite uniform dimension. In the other hand $Soc(M)$ is an essential submodule of M , so M have finite uniform dimension.

1. $\implies$ 4. Since M is an essential submodule of $E(M)$, it follows that $E(M)$ have finite uniform dimension. Thus $E(M)$ is weakly co-Hopfian.

4. $\implies$ 1. Let $E(M)$ be weakly co-Hopfian. Since R is weakly SCI- ring, then $E(M)$ is finitely cogenerated. M as a submodule of $E(M)$ is also finitely cogenerated. It follows that M have finite uniform dimension.

Recall that a ring R is weakly FGI - ring if any weakly co-Hopfian R - module is finitely generated.

Corollary 3.12. *Let R be a commutative ring. The following conditions are equivalent :*

1. R is a weakly SCI-ring.
2. R is an Artinian principal ideal ring.
3. R is a weakly FGI-ring.

Proof It follows from theorem 3.11. and [4].

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**FRACTIONS CONTINUES DE MATRICES ET DÉVELOPPEMENTS
DE CERTAINES FONCTIONS SPÉCIALES**

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Abstract

Abstract. In this paper we summarize the basic definitions and results on matrices and continued fractions. In the convergence case, the continued fractions expansions have the advantage that they converge more rapidly than other numerical algorithms. Further the aim of this paper is to provide some properties and results of continued fractions with matrix arguments. Then we prove the convergence and give the continued fractions expansions of $\sin(A)$, $\sinh(A)$, $\cos(A)$, $\cosh(A)$ and $\arccos(A)$ where A is a matrix.

MSC : 15A60, 40A15, 47A63.

Keywords : matrix continued fraction, convergence, series expansion.

1. INTRODUCTION

En utilisant les approximants de Padé, équivalents aux convergentes des fractions continues de l'exponentielle d'une matrice, R.S. Varga [Var] a étudié et analysé certaines approximations pour des solutions d'équations différentielles paraboliques, comme application cela a donné un exemple de fractions continues non commutatives.

Ainsi, la théorie des fractions continues, dont les éléments sont des opérateurs ont de nombreuses applications dans divers domaines mathématiques et permettent, une accélération de la convergence de certaines approximations. Dans cette direction, et plus généralement dans un espace de Banach, peu de résultats de convergence sur les fractions continues non commutatives sont connus. Deux théorèmes sont cités par P. Wynn [Wyn], où il donne quelques aspects pour la théorie des fractions continues dont les éléments ne commutent pas.

Dans un espace de Banach, des généralisations de certains résultats du cas réel sont publiés par Haydan [Hay] et Negoescu [Neg].

Dans [Rai,Kac], signalons que nous avons donné un critère qui fournit une condition nécessaire pour la convergence de fraction continue de matrices non commutatives de la forme $K(I/A_n)$, que nous avons généralisé au cas de $K(B_n/A_n)$. Nous y avons aussi exprimé une condition suffisante de convergence de $K(B_n/A_n)$.

Dans ce travail, nous citons au début quelques résultats qui lient les éléments, les numérateurs et les dénominateurs des convergentes des fractions continues de matrices. Puis nous rappelons certaines relations fondamentales de ces fractions dont nous aurons besoin par la suite. Vers la fin, nous donnons les développements en fraction continue de certaines fonctions trigonométriques ainsi que ceux des fonctions hyperboliques d'une matrice A , et ceci après avoir prouvé la convergence de ces développements.

2. DÉFINITIONS ET NOTATIONS

2.1. Généralités

On note par M_m l'algèbre des matrices carrées d'ordre m et à coefficients réels muni de la norme classique

$$\forall A \in M_m, \|A\| = \sup_{x \neq 0} \frac{\|Ax\|}{\|x\|} = \sup\{\|Ax\|, \|x\| = 1\}.$$

Soient A et B deux matrices, avec B inversible, on note $\frac{A}{B} = B^{-1}A$, en particulier $\frac{I}{B} = B^{-1}$ où I désigne la matrice identité de M_m .

Il est clair pour toute matrice inversible C , on a

$$\frac{CA}{CB} = \frac{A}{B} \neq \frac{AC}{BC} \text{ et que } C \frac{A}{B} D = \frac{AD}{BC^{-1}} \text{ pour toute matrice } D \in M_m.$$

On dit qu'une matrice A est semidéfinie positive (resp. définie positive) si :

$$\forall x \in \mathbb{R}^m \quad \langle Ax, x \rangle \geq 0 \quad (\text{resp. } \forall x \in \mathbb{R}^m \quad x \neq 0, \langle Ax, x \rangle > 0)$$

où $\langle \cdot, \cdot \rangle$ est le produit scalaire usuel sur $\mathbb{R}^m$ défini par

$$x = (x_1, \dots, x_m) \in \mathbb{R}^m, \quad y = (y_1, \dots, y_m) \in \mathbb{R}^m : \langle x, y \rangle = \sum_{i=1}^m x_i y_i.$$

Ainsi si A et B sont deux matrices symétriques, on note $A \leq B$ si la matrice $B - A$ est semidéfinie positive.

On peut vérifier facilement que si $A \leq B$ alors $AC \leq BC$ pour toute matrice symétrique C .

Définition : Soient $(A_n)_{n \geq 0}, (B_n)_{n \geq 1}$ deux suites de M_m , on appelle fraction continue qu'on note par $K(B_n/A_n)$ la quantité

$$A_0 + \frac{B_1}{A_1 + \dots + B_n/A_n} = \left[A_0; \frac{B_1}{A_1}, \dots, \frac{B_n}{A_n} \right].$$

La matrice $\frac{B_k}{A_k}$ est appelée le terme de rang k ou le $k^{\text{ième}}$ quotient partiel de la fraction continue $K(B_n/A_n)$.

La matrice $\left[A_0; \frac{B_k}{A_k}\right]_{k=1}^n$ s'appelle la $n^{\text{ième}}$ convergente de $K(B_n/A_n)$ et on la note par $Q_n^{-1}P_n$.

Quand $B_n = I$ pour tout $n \geq 1$, alors $K(I/A_n)$ est dite fraction continue ordinaire.

Proposition 1 : Pour $n \geq -1$, les éléments P_n et Q_n de la $n^{\text{ième}}$ convergente de $K(B_n/A_n)$ sont donnés par les relations de récurrence suivantes

$$(1) \quad \begin{cases} P_{-1} = I, P_0 = A_0 \\ Q_{-1} = 0, Q_0 = I \end{cases} \quad \text{et} \quad \begin{cases} P_n = A_n P_{n-1} + B_n P_{n-2} \\ Q_n = A_n Q_{n-1} + B_n Q_{n-2} \end{cases} \quad n \geq 1.$$

Preuve : On utilise une récurrence sur n .

Proposition 2. Soient C et D deux matrices de M_m . avec C inversible, alors on a

$$C \left[A_0, \frac{B_k}{A_k} \right]_{k=1}^n D = \left[CA_0D; \frac{B_1D}{A_1C^{-1}}; \frac{B_2C^{-1}}{A_2}; \frac{B_k}{A_k} \right]_{k=3}^n.$$

2.2. Définition de fraction continue de matrice convergente

On dit que la fraction continue $K(B_n/A_n)$ converge si la suite $(F_n) = \left(\frac{P_n}{Q_n}\right) = (Q_n^{-1}P_n)$ converge dans M_m dans le sens où il existe une matrice $F \in M_m$ telle que $\lim_{n \rightarrow +\infty} \|F_n - F\| = 0$.

Dans ce cas on écrit $\left[A_0; \frac{B_k}{A_k}\right]_{k=1}^{+\infty}$.

Si la fraction continue ci-dessus est convergente et vérifie le fait suivant :

$\forall k \geq 1, A_k = A$ et $B_k = B$ alors on écrit $\left[A_0; \frac{B}{A}\right]_{k=1}^{+\infty}$.

Il est clair que :

$$(2) \quad \frac{P_n}{Q_n} = A_0 + \sum_{i=1}^n \left(\frac{P_i}{Q_i} - \frac{P_{i-1}}{Q_{i-1}} \right).$$

A partir de (2), on voit que la fraction continue $K(B_n/A_n)$ converge dans M_m si et seulement si la série

$$\sum_{n=1}^{+\infty} \left(\frac{P_n}{Q_n} - \frac{P_{n-1}}{Q_{n-1}} \right) \text{ converge dans } M_m.$$

3. CONVERGENCE DE FRACTION CONTINUE DE LA FORME $K(B_n/A_n)$

3.1. Fractions continues équivalentes

Définition : Soient (A_n) , (B_n) , (C_n) et (D_n) quatre suites de matrices de M_m . Les fractions continues $A_0 + K(B_n/A_n)$ et $C_0 + K(D_n/C_n)$ dont les $n^{\text{ièmes}}$ convergentes sont respectivement F_n et G_n sont dites équivalentes si : $\forall n \geq 1 \quad F_n = G_n$.

Maintenant on va énoncer deux lemmes concernant les fractions continues réelles qui nous seront d'une grande importance dans la suite.

Lemme 1. (voir [Jon, Thr])

1) Soit (r_n) une suite de nombre réels non nuls. Alors les fractions continues

$$\left[a_0; \frac{b_1}{a_1}, \frac{b_2}{a_2}, \frac{b_3}{a_3}, \dots \right] \quad \text{et} \quad \left[a_0; \frac{r_1 b_1}{r_1 a_1}, \frac{r_2 r_1 b_2}{r_2 a_2}, \frac{r_3 r_2 b_3}{r_3 a_3}, \dots \right]$$

sont équivalentes.

2) Les fractions continues $a_0 + K(b_n/a_n)$ et $a_0 + K(\beta_n/1)$ sont équivalentes, où

$$\beta_1 = \frac{b_1}{a_1} \quad \text{et} \quad \beta_n = \frac{b_n}{a_n a_{n-1}}, \quad n \geq 2.$$

A partir du développement d'une fonction en série de Taylor, nous donnons dans le lemme suivant le développement en fractions continues de cette série qui a été établi par Euler.

Lemme 2. (voir [Kho])

Soit f une fonction dont le développement en série de Taylor est

$$f(x) = \sum_{n=0}^{+\infty} c_n x^n$$

dans un domaine $D \subset \mathbb{R}$. Alors le développement en fractions continues de $f(x)$ est

$$\begin{aligned} f(x) &= \sum_{n=0}^{+\infty} c_n x^n = \left[c_0; \frac{c_1 x}{1}, \frac{-c_2 x}{c_1 + c_2 x}, \frac{-c_1 c_3 x}{c_2 + c_3 x}, \dots, \frac{-c_{n-2} c_n x}{c_{n-1} + c_n x}, \dots \right]. \\ &= \left[\frac{c_0}{1}, \frac{-c_1 x}{c_0 + c_1 x}, \frac{-c_0 c_2 x}{c_1 + c_2 x}, \frac{-c_1 c_3 x}{c_2 + c_3 x}, \dots, \frac{-c_{n-2} c_n x}{c_{n-1} + c_n x}, \dots \right]. \end{aligned}$$

Remarque : En calculant les premières convergentes de la fraction continue $A_0 + K(B_n/A_n)$, on obtient

$$F_1 = A_0 + A_1^{-1} B_1 = A_0 + (B_1^{-1} A_1)^{-1}$$

$$F_2 = A_0 + (A_1 + A_2^{-1} B_2)^{-1} B_1$$

$$= A_0 + (B_1^{-1} A_1 + (B_2^{-1} A_2 B_1^{-1})^{-1})^{-1}.$$

Si on pose, $A_1^* = B_1^{-1} A_1$ et $A_2^* = B_2^{-1} A_2 B_1$, on a $F_1 = A_0 + \frac{I}{A_1^*}$,

$F_2 = A_0 + \frac{I}{A_1^* + A_2^*}$ et d'une façon générale on prouve par récurrence que si on pose pour tout $k \geq 1$

$$A_{2k}^* = (B_{2k} \dots B_2)^{-1} A_{2k} B_{2k-1} \dots B_1$$

et

$$A_{2k+1}^* = (B_{2k+1} \dots B_1)^{-1} A_{2k+1} B_{2k} \dots B_2,$$

alors les fractions continues $A_0 + K(B_n/A_n)$ et $A_0 + K(I/A_n^*)$ sont équivalentes.

3.2. Critères de convergence de $K(B_n/A_n)$

Dans le théorème suivant, nous énonçons une condition nécessaire pour qu'une fraction continue $K(I/A_n)$ converge dans M_m .

Théorème 1. (voir [Rai,Kac]) Soit $(A_n)_{n \geq 0}$ une suite de M_m . Si la fraction continue $K(I/A_n)$ converge dans M_m alors la série $\sum_{n=0}^{\infty} \|A_n\|$ diverge.

A partir de ce théorème, nous déduisons le corollaire suivant :

Corollaire 1. Si la série $\sum_{n=0}^{\infty} \|A_n\|$ converge absolument dans M_m , alors la fraction continue $K(I/A_n)$ diverge dans M_m .

Nous rappelons ici un critère qui donne une condition suffisante de convergence de fraction continue dans M_m .

Théorème 2. (Voir [Rai,Kac]) Soient $(A_n), (B_n)$ deux suites de M_m . Si

$$\|(B_{2k-2} \dots B_2)^{-1} A_{2k-1}^{-1} B_{2k-1} \dots B_1\| \leq \alpha \text{ et}$$

$$\|(B_{2k-1} \dots B_1)^{-1} A_{2k}^{-1} B_{2k} \dots B_2\| \leq \beta,$$

pour tout $k \geq 1$, où $0 < \alpha < 1$, $0 < \beta < 1$ et $\alpha\beta \leq 1/4$ alors la fraction continue $K(B_n/A_n)$ converge dans M_m .

Remarques : i) Nous remarquons que nous pouvons remplacer la norme utilisée auparavant par n'importe quelle autre norme matricielle.

ii) Signalons aussi que dans [Kac2] nous avons donné le développement en fraction continue de $\tanh(A)$ ainsi que celui d'une solution d'une équation algébrique de Riccati, à savoir

$$(3) \quad XCX = XA + A^T X + B,$$

où B et C sont deux matrices carrées symétriques de M_m et définies positives et A une matrice quelconque de M_m .

4. RÉSULTAT PRINCIPAL

4.1. Convergence de $\sin(A)$ et de $\sinh(A)$

Proposition 3. Soit A une matrice de M_m telle que $\|A\| = \alpha$ où α est un réel $< \frac{1}{2}$. La fraction continue

$$(4) \quad \left[0; \frac{A}{I}, \frac{A^2}{2.3I - A^2}, \frac{(2n-4)(2n-3)A^2}{(2n-2)(2n-1)I - A^2} \right]_{n=3}^{+\infty}$$

converge dans M_m . De plus, le développement en fraction continue ci-dessus est celui de $\sin(A)$ c'est à dire que

$$\sin(A) = \lim_{n \rightarrow +\infty} \left[0; \frac{A}{I}, \frac{A^2}{2.3I - A^2}, \frac{2.3A^2}{4.5I - A^2}, \dots, \frac{(2n-4)(2n-3)A^2}{(2n-2)(2n-1)I - A^2} \right]$$

On a besoin du lemme suivant.

Lemme 3. Soit x un nombre réel non nul et dont le carré n'est pas entier positif. Le développement en fraction continue de $\sin(x)$ est donné par

$$\sin x = \lim_{n \rightarrow +\infty} \left[0; \frac{x}{1}, \frac{x^2}{2.3 - x^2}, \frac{2.3x^2}{4.5 - x^2}, \dots, \frac{(2n-4)(2n-3)x^2}{(2n-2)(2n-1) - x^2} \right].$$

Preuve. A partir du lemme 2 et en utilisant certaines transformations données par la proposition 2, on déduit le développement en fraction continue de la fonction $\sin(x)$ où x est un réel non nul dont le carré n'est pas un entier positif.

Preuve de la proposition 3. Soit A une matrice de M_m telle que $\|A\| = \alpha < 1/2$, nous étudions la convergence de la fraction continue $A_0 + K(B_k/A_k)$ avec $A_0 = 0$,

$$\begin{cases} A_1 = I \\ B_1 = A, B_2 = A^2 \end{cases}, \begin{cases} A_k = (2k-2)(2k-1)I - A^2, k \geq 2 \\ B_k = (2k-4)(2k-3)A^2, k \geq 3. \end{cases}$$

Vérifions que les conditions du théorème 2 sont vérifiées. Pour $k \geq 1$, nous avons

$$(B_{2k-2} \dots B_2)^{-1} A_{2k-1}^{-1} B_{2k-1} \dots B_1 = (4k-6)(4k-5) \times \frac{(4k-10)(4k-9)}{(4k-8)(4k-7)} \dots \frac{(2.3)}{(4.5)} \left(((4k-4)(4k-3)I - A^2) A^{2(k-1)} \right)^{-1} A^{2k-1}.$$

Or les deux matrices $(4k-4)(4k-3)I - A^2$ et A^{2k-1} commutent, donc l'égalité ci-dessus devient

$$(B_{2k-2}...B_2)^{-1} A_{2k-1}^{-1} B_{2k-1}...B_1 = (4k-6)(4k-5) \times \frac{(4k-10)(4k-9)}{(4k-8)(4k-7)} \dots \frac{(2.3)}{(4.5)} \left((4k-4)(4k-3)I - A^2 \right)^{-1} A.$$

De plus, puisque $\|A\| < 1$, alors d'une part, on a

$$\left\| \left((4k-4)(4k-3)I - A^2 \right)^{-1} \right\| < \frac{1}{(4k-4)(4k-3)} \times \left(1 - \frac{\|A\|^2}{(4k-4)(4k-3)} \right)^{-1},$$

et d'autre part pour k assez grand, on a

$$\left(1 - \frac{\|A\|^2}{(4k-4)(4k-3)} \right)^{-1} < 1.$$

Ce qui implique que

$$\|(B_{2k-2}...B_2)^{-1} A_{2k-1}^{-1} (B_{2k-1}...B_1)\| < \frac{(4k-6)(4k-5)}{(4k-4)(4k-3)} \|A\| < 1/2.$$

Il nous reste à prouver la deuxième inégalité du théorème 2, on a

$$(B_{2k-1}...B_1)^{-1} A_{2k}^{-1} B_{2k}...B_2 = (4k-4)(4k-3) \frac{(4k-8)(4k-7)}{(4k-6)(4k-5)} \dots \frac{1}{(2.3)} \times \left(((4k-4)(4k-3)I - A^2) A^{2k-1} \right)^{-1} A^{2k}.$$

En utilisant encore le fait que les matrices $((4k-2)(4k-1)I - A^2)$ et $A^{2(k-1)}$ commutent et en passant à la norme, on trouve

$$\|(B_{2k-1}...B_1)^{-1} A_{2k}^{-1} (B_{2k}...B_2)\| < \frac{(4k-4)(4k-3)}{(4k-2)(4k-1)} \times \left\| \left(I - \frac{A^2}{(4k-4)(4k-3)} \right)^{-1} \right\| \|A\| < \|A\| = \alpha < 1/2,$$

pour k assez grand. Ce qui achève la preuve de la proposition 3.

Corollaire 2. Soit A une matrice de M_m telle que $\|A\| = \alpha$ où α est un réel $< \frac{1}{2}$. Le développement en fraction continue de la fonction sinus hyperbolique $\sinh(A)$ est

$$\sinh(A) = \lim_{n \rightarrow +\infty} \left[0; \frac{A}{I}, \frac{-A^2}{2.3I + A^2}, \frac{-2.3A^2}{4.5I + A^2}, \dots, \frac{-(2n-4)(2n-3)A^2}{(2n-2)(2n-1)I + A^2} \right]$$

Pour la preuve du corollaire 2, le développement en fraction continue de $\sinh(A)$ s'obtient en utilisant la relation suivante,

$$\sinh(A) = -i \sin(iA).$$

4.2. Convergence de $\cos(A)$ et $\cosh(A)$

Proposition 4. Soit A une matrice de M_m telle que $\|A\| = \alpha$ où α est un réel $< \frac{1}{2}$. La fraction continue

$$(5) \quad \left[0; \frac{I}{I}, \frac{A^2}{1.2I - A^2}, \frac{(2n-5)(2n-4)A^2}{(2n-3)(2n-2)I - A^2} \right]_{n=3}^{+\infty}$$

converge dans M_m . De plus, le développement en fraction continue ci-dessus est celui de $\cos(A)$ c'est à dire que

$$\cos(A) = \lim_{n \rightarrow +\infty} \left[0; \frac{I}{I}, \frac{A^2}{1.2I - A^2}, \frac{2.3A^2}{4.5I - A^2}, \dots, \frac{(2n-4)(2n-3)A^2}{(2n-2)(2n-1)I - A^2} \right]$$

Puisque $\cosh(A) = \cos(iA)$, en remplaçant la matrice A par iA dans la proposition 4, on obtient le corollaire suivant.

Corollaire 3. Soit A une matrice de M_m telle que $\|A\| = \alpha$ où α est un réel $< \frac{1}{2}$. Le développement en fraction continue de la fonction sinus hyperbolique $\cosh(A)$ est

$$\cosh(A) = \lim_{n \rightarrow +\infty} \left[0; \frac{I}{I}, \frac{-A^2}{2.3I + A^2}, \frac{-2.3I}{4.5I + A^2}, \dots, \frac{-(2n-4)(2n-3)A^2}{(2n-2)(2n-1)I + A^2} \right]$$

Preuve de la proposition 4. En se basant sur les résultats des lemmes 1 et 2, on prouve que pour tout x dont le carré n'est pas un entier positif que

$$\cos x = \lim_{n \rightarrow +\infty} \left[0; \frac{1}{1}, \frac{x^2}{1.2 - x^2}, \frac{2.3x^2}{4.5 - x^2}, \dots, \frac{(2n-4)(2n-3)x^2}{(2n-2)(2n-1) - x^2} \right].$$

Soit A une matrice de M_m telle que $\|A\| = \alpha < 1/2$, on prouve de la même manière que dans 4.1) que la fraction continue $A_0 + K(B_k/A_k)$ avec $A_0 = 0$,

$$\begin{cases} A_1 = I \\ B_1 = I, B_2 = A^2 \end{cases} \begin{cases} A_k = (2k-3)(2k-2)I - A^2I, k \geq 2 \\ B_k = (2k-5)(2k-4)A^2, k \geq 3 \end{cases}$$

converge pour toute matrice $A \in M_m$ telle que $\|A\| = \alpha < \frac{1}{2}$.

4.3. Convergence de $\arccos(A)$

Proposition 5. Soit A une matrice de M_m telle que $\|A\| < \frac{1}{2}$. La fraction continue

$$\left[0; \frac{A}{(I - A^2)^{1/2}}, \frac{-b_n(I - A^2)}{I} \right]_{n=2}^{+\infty}$$

avec

$$b_{2n} = \frac{2n(2n-1)}{(4n-1)(4n-3)} \text{ et } b_{2n+1} = \frac{2n(2n-1)}{(4n+1)(4n-1)}, n \geq 1,$$

converge dans M_m . De plus, le développement en fraction continue ci-dessus est celui de $\arccos(A)$, c'est à dire que

$$\arccos(A) = \lim_{n \rightarrow +\infty} \left[0; \frac{A}{(I - A^2)^{1/2}}, \frac{-b_n(I - A^2)}{I} \right]_{n=2}^{+\infty}.$$

Preuve de la proposition 5. Signalons tout d'abord que la suite des termes $(b_k)_{k \geq 2}$ est croissante. Le lemme suivant nous sera d'une grande importance pour cette preuve.

Lemme 4 (Voir [Han]). Soit x un nombre réel tel que $|x| < 1$, le développement en fraction continue de $\arccos x$ est donné par :

$$\arccos x = \lim_{n \rightarrow +\infty} \left[0; \frac{x\sqrt{1-x^2}}{1}, \frac{-b_n\sqrt{1-x^2}}{1} \right]_{n=2}^{+\infty} \text{ avec}$$

$$b_{2n} = \frac{2n(2n-1)}{(4n-1)(4n-3)} \text{ et } b_{2n+1} = \frac{2n(2n-1)}{(4n+1)(4n-1)}, n \geq 1.$$

Soit A une matrice de M_m telle que $\|A\| < 1$, nous étudions la convergence de la fraction continue $A_0 + K(B_k/A_k)$ avec $A_0 = 0$,

$$\begin{cases} A_1 = I \\ B_1 = A(I - A^2)^{1/2} \end{cases} \quad \begin{cases} B_{2n} = \frac{2n(2n-1)}{(4n-1)(4n-3)}(I - A^2), n \geq 1 \\ B_{2n+1} = \frac{2n(2n-1)}{(4n+1)(4n-1)}(I - A^2), n \geq 1. \end{cases}$$

En tenant compte de la croissance des b_k définis dans la proposition 5, on obtient

$$\|(B_{2k-1} \dots B_3)^{-1} A_{2k}^{-1} (B_{2k} \dots B_4)\| < \frac{2k(2k-1)}{(4k-1)(4k-3)}.$$

Comme on s'intéresse à la convergence de la fraction continue $A_0 + K(B_k/A_k)$, donc k est assez grand, on remarque que la première inégalité du théorème 2 est vérifiée.

On prouve de la même manière la deuxième inégalité du théorème 2. D'où le résultat de la proposition 5.

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QUASI δ -DISCRETE MODULES**Mohammad Javad Nematollahi**

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Abstract

A module M is called quasi δ -discrete if M is δ -lifting and M has (D_3) . We prove that if M is a quasi δ -discrete module such that every δ -supplement submodule is δ -coclosed, then M has a decomposition $M = \oplus_{i \in I} H_i$ where each H_i is hollow. Moreover, if M is amply supplemented, such a decomposition complements direct summands, and hence is unique up to isomorphism. Also we show that such a module M has a decomposition unique up to isomorphism, $M = M_1 \oplus M_2$, where $\delta(M_1) \ll_{\delta} M_1$ and $\delta(M_2) = M_2$.

2000 Mathematics Subject Classification: 16L30, 16E50**Keywords:** δ -supplemented module, δ -lifting module, quasi δ -discrete module

1 Introduction

Throughout this article, all rings are associative with an identity, and all modules are unitary right R -modules. A submodule L of a module M is called *small* in M (denoted by $L \ll M$), if for every proper submodule K of M , $L + K \neq M$. A module M is called *hollow*, if every proper submodule of M is small in M .

For two submodules N and K of M , N is called a *supplement* of K in M if, N is minimal with the property $M = K + N$, equivalently $M = K + N$ and $N \cap L \ll N$. A module M is called *supplemented* if, every submodule of M has a *supplement* in M . Also M is called *amply supplemented* if, for any two submodules L, K of M with $M = L + K$ there exists a *supplement* P of L such that $P \leq K$. A module M is called *weakly supplemented* if, for each submodule A of M there exists a submodule B of M such that $M = A + B$ and $A \cap B \ll M$, in this case B is called a *weak supplement* of A in M . Let M be a module and $B \leq A \leq M$. If $A/B \ll M/B$, then B is called a *cosmall* submodule of A in M . A submodule A of M is called *coclosed* in M if, A has no proper cosmall submodule. The relationship between supplements and coclosed submodules is given in the next Lemma:

Lemma 1.1 [1, 20.2] *Let N be a submodule of a module M . Consider the following statements:*

- (1) N is a supplement in M ;
- (2) N is coclosed in M ;
- (3) For all $K \subseteq N$, $K \ll M$ implies $K \ll N$.

Then (1) $\Rightarrow$ (2) $\Rightarrow$ (3) hold and if M is weakly supplemented then (3) $\Rightarrow$ (1) holds.

A submodule K of M is called *essential* in M (denoted by $K \leq^{ess} M$) if, $K \cap X \neq 0$ for every non zero submodule X of M . The *singular submodule* of a module M (denoted by $Z(M)$) is $Z(M) = \{x \in M \mid Ix = 0; \text{ for some ideal } I \leq^{ess} R\}$. A module M is called *singular* (*nonsingular*) if, $Z(M) = M$ ($Z(M) = 0$).

Let M be a module. A submodule N of M is said to be δ -small in M (notation $N \ll_{\delta} M$) if, whenever $N + X = M$ with M/X singular, we have $X = M$. The concept of δ -small submodules was introduced by Zhou in [5]. A module M is called δ -hollow, if every proper submodule of M is δ -small in M .

Every small submodule of M is δ -small in M and the converse is true when M is singular. But as we see in the next example the converse need not be true in general:

Example 1.2 *Let R be a right semisimple ring and M be a nonzero right R -module. Then M is nonsingular and semisimple. For any nonzero $N \leq M$, N is a direct summand of M and hence is not small in M , but every submodule of M (even M itself) is δ -small in M .*

Let N and L be submodules of a module M . N is called a (weak) δ -supplement of L in M , if $N + L = M$ and $N \cap L \ll_{\delta} N$ ($N \cap L \ll_{\delta} M$). A module M is called (weakly) δ -supplemented if every submodule of M has a (weak) δ -supplement in M . M is called *amply δ -supplemented* if for any submodules A and B of M with $M = A + B$, A has a δ -supplement contained in B .

Now a question arises as follow:

Does Lemma 1.1 hold for δ -supplement and δ -coclosed submodules?

In general the answer of this question is negative. Indeed Example 1.2 shows that even in an amply supplemented module it does not need a δ -supplement submodule (and even a direct summand) to be δ -coclosed. In this example any nontrivial submodule of M is a direct summand (and hence a (δ -)supplement), but it is not δ -coclosed in M .

We recall that a family $\{X_{\lambda} : \lambda \in \Lambda\}$ of submodules of a module M is called a *local direct summand* of M , if $\sum_{\lambda \in \Lambda} X_{\lambda}$ is direct and $\sum_{\lambda \in F} X_{\lambda}$ is a direct summand of M for every finite subset $F \subseteq \Lambda$. If even $\sum_{\lambda \in \Lambda} X_{\lambda}$ is a direct summand of M , we say that the local direct summand is a direct summand. It is well known that if every local direct summand of a module M is a direct summand, then M is a direct sum of indecomposable modules (see for example [3, Theorem 2.17]).

We recall that M is called a (D_3) -module (or we say M has (D_3)) if for any two direct summands A and B of M with $M = A + B$, $A \cap B$ is a direct summand of M .

We denote by $Rad(M)$ the radical of M which is the sum of all small submodules of M . Also we write $A \leq^m M$ to indicate that A is a maximal submodule of M and We write $A \leq^d M$ to indicate that A is a direct summand of M .

In Sections 2 and 3, various properties of δ -supplemented and δ -lifting modules are obtained. In Section 4, we introduce quasi δ -discrete modules and give some decomposition theorems about them.

2 δ -Supplemented Modules

Lemma 2.1 [5, Lemma 1.3] *Let M be a module.*

- (1) *For submodules N, K and L of M with $K \subseteq N$ we have:*
 - (a) *$N \ll_\delta M$ if and only if $K \ll_\delta M$ and $N/K \ll_\delta M/K$.*
 - (b) *$N + L \ll_\delta M$ if and only if $N \ll_\delta M$ and $L \ll_\delta M$.*
- (2) *If $K \ll_\delta M$ and $f : M \rightarrow N$ is a homomorphism, then $f(K) \ll_\delta N$. In particular, if $K \ll_\delta M \subseteq N$, then $K \ll_\delta N$.*
- (3) *Let $K_1 \leq M_1 \leq M$, $K_2 \leq M_2 \leq M$, and $M = M_1 \oplus M_2$. Then $K_1 \oplus K_2 \ll_\delta M_1 \oplus M_2$ if and only if $K_1 \ll_\delta M_1$ and $K_2 \ll_\delta M_2$.*

Lemma 2.2 *Let N and L be submodules of a module M . Then the following are equivalent*

- (1) *N is a δ -supplement of L in M ;*
- (2) *$N + L = M$ and for each $K \leq N$ with $K + L = M$ and N/K singular, $K = N$.*

Proof The proof is straightforward. □

Proposition 2.3 *Let M be a module and $N \leq M$. Consider the following conditions:*

- (1) *N is a δ -supplement submodule of M ;*

(2) N is weak δ -coclosed in M ;

(3) For all $X \leq M$, $X \ll_\delta M$ implies $X \ll_\delta N$.

Then (1) $\Rightarrow$ (2) $\Rightarrow$ (3) hold. If M is weakly δ -supplemented, then (3) $\Rightarrow$ (1) holds.

Proof (1) $\Rightarrow$ (2) Let N be a δ -supplement of a submodule K in M . Then $N + K = M$ and $N \cap K \ll_\delta N$. Suppose that for $X \leq N$, N/X is singular and $N/X \ll_\delta M/X$. Then we have $N/X + (K + X)/X = M/X$ and since $M/(K + X)$ is a homomorphic image of the singular module N/X is singular. Hence since $N/X \ll_\delta M/X$, we have $K + X = M$. Now by Lemma 2.2 we get $X = N$, i. e. N is δ -coclosed in M .

(2) $\Rightarrow$ (3) Let $K \ll_\delta M$ and let $K + X = N$ for some $X \leq N$ with N/X singular. Then $N/X = (K + X)/X$ is the image of a δ -small submodule and hence by Lemma 2.1(2), $N/X \ll_\delta M/X$. So we have $X = N$, because N is δ -coclosed.

(3) $\Rightarrow$ (1) Let L be a weak δ -supplement of N in M . Then $N + L = M$ and $N \cap L \ll_\delta M$. Then $N \cap L \ll_\delta N$, by (3). It means that N is a δ -supplement of L in M . $\square$

Lemma 2.4 Let M be a module. Then for any $a \in M$ we have: aR is not δ -small in M , if and only if there exists a maximal submodule C of M with M/C singular and $a \notin C$.

Proof The proof is straightforward by a standard argument and using Zorn's Lemma. $\square$

The next Definition is essentially due to Zhou [5].

Definition 2.5 Let $\wp$ be the class of all singular simple modules. For a module M let $\delta(M) = \text{Rej}_M(\wp) = \cap \{N \subseteq M \mid M/N \in \wp\}$ be the reject of $\wp$ in M .

From the definition we immediately have for any module M , $\delta(M/\delta(M)) = 0$.

Proposition 2.6 Given a module M , each of the following sets is equal to $\delta(M)$:

- (1) $A_1 = \sum \{A \mid A \ll_\delta M\}$;
- (2) $A_2 = \cap \{B \mid B \leq^M M \text{ with } M/B \text{ singular}\}$;
- (3) $A_3 = \cap \{\ker \phi \mid \phi \in \text{Hom}(M, N) \text{ such that } N \text{ is singular simple}\}$;
- (4) $A_4 = \cap \{\ker \phi \mid \phi \in \text{Hom}(M, N) \text{ such that } N \text{ is singular semisimple}\}$.

Proof A_3 is just the definition of $\delta(M)$. Clearly $A_4 \subseteq A_3$ and it is easy to see that $A_3 \subseteq A_2$. Using Lemma 2.4 shows that $A_2 \subseteq A_1$. To show that $A_1 \subseteq A_4$ let $A \ll_\delta M$ and $\phi : M \rightarrow N$ be a homomorphism such that N is singular semisimple. By [5, Lemma 1.3] $\phi(A) \ll_\delta N$. Since N is singular we have $\phi(A) \ll N$. So $\phi(A) = 0$ because N is semisimple. Hence $A \subseteq \ker \phi$. Therefore $A_1 \subseteq A_4$. $\square$

Proposition 2.7 *Let U and V be submodules of a module M . Assume V to be a δ -supplement of U in M . Then*

- (1) *If $W + V = M$ for some $W \subseteq U$, then V is a δ -supplement of W in M ;*
- (2) *If $K \ll_\delta M$, then V is a δ -supplement of $U + K$ in M ;*
- (3) *For $K \ll_\delta M$ we have $K \cap V \ll_\delta V$ and so $\delta(V) = V \cap \delta(M)$;*
- (4) *For $L \subseteq U$, $(V + L)/L$ is a δ -supplement of U/L in M/L ;*
- (5) *If $\delta(M) \ll_\delta M$, or $\delta(m) \subseteq U$ and if $p : M \rightarrow M/\delta(M)$ is the canonical projection, then $M/\delta(M) = Up \oplus Vp$.*

Proof (1) follows directly from definition.

(2) If $K \ll_\delta M$, then, for $X \subseteq V$ with $U + K + X = M$ and V/X singular, we have $M/(U + X)$ is a homomorphic image of V/X and hence singular. So $U + X = M$ and hence by Lemma 2.2 we have $V = X$.

(3) Let $K \ll_\delta M$ and $X \subseteq V$ with $(K \cap V) + X = V$ and V/X singular. Then $M = U + V = U + (K \cap V) + X = U + X$ and hence by Lemma 2.2 $X = V$, i.e. $K \cap V \ll_\delta V$. This yields $V \cap \delta(M) \subseteq \delta(V)$. Since $\delta(V) \subseteq V \cap \delta(M)$ always holds, we get the desired equality.

(4) For $L \subseteq U$, we have $U \cap (V + L) = U \cap V + L$ (modular law) and $U/L \cap (V + L)/L = ((U \cap V) + L)/L$. Since $U \cap V \ll_\delta V$, it follows that $((U \cap V) + L)/L \ll_\delta (V + L)/L$ (image of a δ -small submodule). Now the assertion follows from $U/L + (V + L)/L = M/L$.

(5) If $\delta(M) \subseteq U$ then by (4), $Up \cap Vp \ll_{\delta} Vp$. Hence also $Up \cap Vp \ll_{\delta} M/\delta(M)$ and therefore $Up \cap Vp = 0$. If $\delta(M) \ll_{\delta} M$, then by (2) V is a δ -supplement of $U + \delta(M)$. $\square$

Proposition 2.8 *Let M be an amply δ -supplemented module. Then*

- (1) *Every supplement submodule of M is an amply δ -supplemented module;*
- (2) *Direct summands and factor modules of M are amply δ -supplemented.*

Proof The proof is easy and omitted. $\square$

3 δ -Lifting Modules

Proposition 3.1 *Let M be an amply δ -supplemented module. Then every submodule N of M contains a δ -supplement submodule N' such that $N/N' \ll_{\delta} M/N'$.*

Proof If $N \ll_{\delta} M$ then the desired δ -supplement submodule is 0. Let N be a non δ -small submodule of M . Then N has a δ -supplement $K \neq M$. Hence $N + K = M$ and $N \cap K \ll_{\delta} K$. Since M is amply δ -supplemented, there is a submodule N' of N such that $K + N' = M$ and $K \cap N' \ll_{\delta} N'$. We want to show that $N/N' \ll_{\delta} M/N'$. Let L be a submodule of M such that $N' \subseteq L$ and $N/N' + L/N' = M/N'$ with $\frac{M/N'}{L/N'}$ singular. So $N + L = M$ and M/L is singular. By modular law from $K + N' = M$ we have $N = N' + (K \cap N)$ and hence $N' + (K \cap N) + L = M$. So $N' + L = M$ and it implies that $L = M$. $\square$

Proposition 3.2 *For a submodule $U \subseteq M$, the following are equivalent:*

- (1) *There is a direct summand X of M with $X \subseteq U$ and $U/X \ll_{\delta} M/X$;*
- (2) *There is a direct summand $X \subseteq M$ and a submodule Y of M with $X \subseteq U$, $U = X + Y$ and $Y \ll_{\delta} M$;*

- (3) There is a decomposition $M = X \oplus X'$ with $X \subseteq U$ and $X' \cap U \ll_\delta X'$;
- (4) U has a δ -supplement V in M such that $U \cap V$ is a direct summand in U ;
- (5) There is an idempotent $e \in \text{End}(M)$ with $Me \subseteq U$ and $U(1 - e) \ll_\delta M(1 - e)$.

Proof (1) $\implies$ (2) If $M = X \oplus X'$ and $U/X \ll_\delta M/X$, then $U = X + (U \cap X')$. Let $\phi : M/X \rightarrow X'$ be the obvious isomorphism. Then, since $U/X \ll_\delta M/X$, we have $(U/X)\phi = U \cap X' \ll_\delta X'$. Hence $U \cap X' \ll_\delta M$.

(2) $\implies$ (3) If $M = X \oplus X'$, then X' is a δ -supplement of X in M and hence a δ -supplement of $U = X + Y$ in M (see Proposition 2.7(2)), and so $U \cap X' \ll_\delta X'$.

(3) $\implies$ (4) From (3), $U = X \oplus (U \cap X')$ and X' is a δ -supplement of U in M .

(4) $\implies$ (5) By (4), $M = U + V$, $U \cap V \ll_\delta V$ and $U = (U \cap V) \oplus X$ for some $X \subseteq U$. Then $M = U + V = X + (U \cap V) + v = X + V$ and $X \cap V = 0$, and so $M = X \oplus V$. Let $e : M \rightarrow X$ be the projection map along V . Then $Me \subseteq U$. As $X \subseteq U$, we have $U(1 - e) = U \cap M(1 - e) = U \cap V \ll_\delta V = M(1 - e)$, as required.

(5) $\implies$ (1) Letting $X = Me$, we have $M = X \oplus M(1 - e)$ and $X \subseteq U$. Taking $\phi : M/X \rightarrow M(1 - e)$ to be the obvious isomorphism gives $(U/X)\phi = U(1 - e) \ll_\delta M(1 - e)$. As ϕ is an isomorphism, it follows that $U/X \ll_\delta M/X$. $\square$

Definition 3.3 A module M is called δ -lifting if for any $A \leq M$, there exists a decomposition $M = M_1 \oplus M_2$ such that $M_1 \leq A$ and $A/M_1 \ll_\delta M/M_1$.

For example every lifting module is δ -lifting. Also every δ -hollow module is δ -lifting and it is easy to see that every indecomposable δ -lifting module is δ -hollow.

Lemma 3.4 [5, Lemma 1.2] Let N be submodule of M . Then the following are equivalent:

- (1) $N \ll_\delta M$;

- (2) If $X + N = M$, then $M = X \oplus Y$ for a projective semisimple submodule Y with $Y \subseteq N$.

Proposition 3.5 *For an indecomposable module M the following are equivalent:*

- (1) M is δ -lifting;
- (2) M is δ -hollow;
- (3) M is hollow;
- (4) M is lifting.

Proof (1) $\implies$ (2), (3) $\implies$ (4) and (4) $\implies$ (1) are obvious. Now (2) $\implies$ (3) follows from Lemma 3.4. $\square$

The next Proposition immediately follows from Proposition 3.2 and also can be found in [2, Lemma 2.3]:

Proposition 3.6 *For a module M the following are equivalent:*

- (1) M is δ -lifting;
- (2) For every submodule N of M there is a decomposition $M = M_1 \oplus M_2$ such that $M_1 \subseteq N$ and $N \cap M_2 \ll_{\delta} M$;
- (3) Every submodule N of M can be written as $N = N_1 \oplus N_2$ with N_1 a direct summand of M and $N_2 \ll_{\delta} M$.

Corollary 3.7 *Every direct summand of a δ -lifting module is δ -lifting.*

Proposition 3.8 *Let M be a δ -lifting module. Then*

- (1) Any δ -coclosed submodule of M is a direct summand;
- (2) M is amply δ -supplemented;
- (3) If $N \subseteq M$ is a fully invariant submodule of M , then M/N is a δ -lifting module.

Proof The proofs are similar to those for lifting case ([1, 22.2]) and omitted. $\square$

Proposition 3.9 *Let M be an amply δ -supplemented module such that every δ -supplement submodule of M is a direct summand. Then M is δ -lifting.*

Proof Let $N \leq M$. If $N \ll_{\delta} M$, then $N/0 \ll_{\delta} M/0$. If N is not a δ -small submodule of M , then by Proposition 3.1 there exists a δ -supplement submodule K of M such that $N/K \ll_{\delta} M/K$. By assumption K is a direct summand of M and hence M is δ -lifting. $\square$

Proposition 3.10 *Let M be a module such that every δ -supplement submodule of M is δ -coclosed in M . Then M is δ -lifting if and only if M is amply δ -supplemented and every δ -supplement submodule is a direct summand.*

Proof The proof can be obtained by Propositions 3.8 and 3.9. $\square$

4 Quasi δ -discrete Modules

Definition 4.1 *A module M is called quasi δ -discrete if M is δ -lifting and M has (D_3) .*

It is clear that every quasi discrete (e. g. hollow) module is quasi δ -discrete and every direct summand of a quasi δ -discrete module is quasi δ -discrete.

Remark 4.2 *As we saw in the Example 1.2 a nonzero module M can be δ -small in itself (unlike the small case). But if every δ -supplement submodule of M is δ -coclosed in M then it is not hard to see that for every $N \leq M$, $N \ll_{\delta} N$ if and only if $N = 0$.*

Remark 4.3 *Again seeing the Example 1.2 shows that there exists a decomposable δ -hollow module (unlike the hollow case). But it is easy to see that if M is a δ -hollow module such that every δ -supplement submodule of M is δ -coclosed in M then M is indecomposable. In this case M is quasi δ -discrete.*

Proposition 4.4 *Let M be a module such that every δ -supplement submodule is δ -coclosed. Then M is quasi δ -discrete if and only if M is amply δ -supplemented, and $M = X \oplus Y$ for any two submodules X and Y of M which are δ -supplement of each other.*

Proof $(\Rightarrow)$ M is amply δ -supplemented and X and Y are direct summands of M by Proposition 3.10. Hence $X \cap Y$ is a direct summand of M by (D_3) . Since $X \cap Y \ll_\delta M$ we have $X \cap Y \ll_\delta X \cap Y$ and hence by Remark 4.2 $X \cap Y = 0$.

$(\Leftarrow)$ By Proposition 3.10 M is δ -lifting. Let A and B be direct summands of M such that $M = A + B$. Now B contains a δ -supplement B' of A which is a direct summand of M . Thus $A \cap B' \ll_\delta B' \leq M$, consequently $A \cap B' \ll_\delta A$. Thus A is a δ -supplement of B' , hence $M = A \oplus B'$ by hypothesis. Then $B = B' \oplus A \cap B$, so $A \cap B$ is a direct summand of M and hence M has (D_3) . $\square$

Lemma 4.5 *Let M be a quasi δ -discrete module such that every δ -supplement submodule of M is δ -coclosed. Let A be a submodule of M , and B a δ -supplement of A . If C is a δ -supplement submodule of M contained in A , then $C \cap B = 0$ and $C \oplus B \leq^d M$.*

Proof B and C are direct summands of M (Proposition 3.10). Writing $M = C \oplus C'$ we get $A = A \cap C' \oplus C$, hence $M = A \cap C' + C + B$. By Proposition 3.10, $A \cap C'$ contains a δ -supplement D of $C + B$, which is a direct summand. Since $D \leq C'$, $D \oplus C \leq^d M$. Now $M = (D \oplus C) + B$, a sum of two summands, so $(D \oplus C) \cap B \leq^d M$ by (D_3) . However $(D \oplus C) \cap B \leq A \cap B \ll_\delta M$. Hence $(D \oplus C) \cap B = 0$, and $M = D \oplus C \oplus B$. $\square$

Corollary 4.6 *Let M be a quasi δ -discrete module such that every δ -supplement submodule is δ -coclosed. Then the union of any chain of direct summands is a direct summand, and local direct summands are direct summands.*

Proof Let A_α be a chain of direct summands of M , and let $A = \bigcup_\alpha A_\alpha$. Let B be a δ -supplement of A . Then $A_\alpha \cap B = 0$ for every

α by Lemma 4.5, and hence $A \cap B = 0$, so $M = A \oplus B$. The last statement follows by [3, Lemma 2.16]. $\square$

Proposition 4.7 *Let M be a quasi δ -discrete module such that every δ -supplement submodule is δ -coclosed. Let N be a direct summand and H be δ -hollow direct summand of M . Then either $N + H \leq^d M$ and $N \cap H = 0$, or $N + H = N \oplus S$ with $S \ll_\delta M$. Moreover if M is amply supplemented then instead the second case we have $N + H = N \oplus S$ with $S \ll M$ and H is isomorphic to a direct summand of N .*

Proof Let L be a δ -supplement of $N + H$, then $N \cap L = 0$ by Lemma 4.5. Now $M = N + L + H$. We consider the two possible cases:

(i) $H \not\leq N \oplus L$: Then since H is δ -hollow, H is a δ -supplement of $N + L$, hence $N \cap H = 0$ and $N \cap H \leq^d M$ by Lemma 4.5.

(ii) $H \leq N \oplus L$: Then $M = N \oplus L$ and $N + H = N \oplus (N + H) \cap L = N \oplus S$, where $S = (N + H) \cap L \ll_\delta M$ since L is a δ -supplement of $N + H$.

For the last statement if M is amply supplemented, let L be a supplement of $N + H$. Then in case (ii) we have $N + H = N \oplus S$, with $S \ll M$. Now write $M = H \oplus H'$. Then $M = N + H + H' = N + S + H' = N + H'$. Let N' be a supplement of H' contained in N . Then $M = N' \oplus H'$ by Proposition 4.4. Hence $H \simeq N' \leq^d N$. $\square$

Lemma 4.8 *Let M be a δ -lifting module such that every δ -supplement submodule is δ -coclosed. Then every direct summand of M has this property.*

Proof The proof is easy and omitted. $\square$

Theorem 4.9 *Let M be a quasi δ -discrete module such that every δ -supplement submodule is δ -coclosed. Then M has a decomposition $M = \bigoplus_{i \in I} H_i$ where each H_i is hollow. Moreover, if M is amply supplemented, such a decomposition complements direct summands, and hence is unique up to isomorphism.*

Proof By Corollary 4.6, $M = \oplus_{i \in I} H_i$ where each H_i is indecomposable. It follows by Proposition 3.5 that each H_i is hollow.

Next we prove that if M is amply supplemented the decomposition $M = \oplus_{i \in I} H_i$ complements direct summands. Let X be any direct summand of M . By Zorn's Lemma and Corollary 4.6, we get a subset $J \subseteq I$ maximal such that $X \cap (\oplus_{j \in J} H_j) = 0$, and $X \oplus (\oplus_{j \in J} H_j) \leq^d M$. Write $M = X \oplus (\oplus_{j \in J} H_j) \oplus Y$. We will complete the proof by showing that $Y = 0$. Assume that $Y \neq 0$. Since Y is quasi δ -discrete and by Lemma 4.8, Y contains a nonzero hollow direct summand H (by what we have already proved). Consequently $M = T \oplus H$ where $X \oplus (\oplus_{j \in J} H_j) \leq T$. If for some $i \in I \setminus J$, $T + H_i \leq^d M$, then $T \cap H_i = 0$, and we get a contradiction to the maximality of J . Thus $T + H_i$ is not a direct summand of M for any $i \in I$. It then follows by Proposition 4.7 that for all $i \in I$, $T + H_i = T \oplus S_i$ with $S_i \ll M$. Consider any finite subset $F \subseteq I$. Then $T + (\oplus_{i \in F} H_i) \neq M$, and $M = T \oplus H = T + (\oplus_{i \in F} H_i) + H$. Hence H is a supplement of $T + (\oplus_{i \in F} H_i)$. But then $H \cap (\oplus_{i \in F} H_i) = 0$ by Lemma 4.5. Since this is true for any finite subset $F \subseteq I$, $H \cap (\oplus_{i \in I} H_i) = 0$, a contradiction. Therefore $Y = 0$, and $M = X \oplus (\oplus_{j \in J} H_j)$. $\square$

Recall that a family of modules $\{M_\alpha : \alpha \in A\}$ is called *locally semi-T-nilpotent*, if for any subfamily $\{M_{\alpha_i} : i \in N\}$ with distinct α_i and any family of non-isomorphisms $\{f_i : M_{\alpha_i} \rightarrow M_{\alpha_{i+1}}\}$, and for every $x \in M_{\alpha_1}$, there exists $n \in N$ (depend on x) such that $f_n \dots f_2 f_1(x) = 0$. It is well known that for a module M any decomposition $M = \oplus_{\alpha \in A} M_\alpha$ which complements direct summand is locally semi-T-nilpotent.

Corollary 4.10 *Let M be an amply supplemented quasi δ -discrete module such that every δ -supplement submodule is δ -coclosed. Then any decomposition of M as a direct sum of hollows is locally semi-T-nilpotent.*

A module M is called local if M has a largest submodule which contains all proper submodules of M . Obviously if M is hollow and $\text{Rad}(M) \neq M$, then M is local.

Proposition 4.11 *Let M be an amply supplemented quasi δ -discrete module such that every δ -supplement submodule is δ -coclosed. Then the following are equivalent:*

- (1) $\delta(M) \ll_{\delta} M$;
- (2) Every proper submodule of M is contained in a maximal submodule;
- (3) M is a direct sum of local modules.

Proof (1) $\Rightarrow$ (3) By Theorem 4.9, $M = \oplus_{i \in I} H_i$ with each H_i hollow. Since $\delta(H_i) \leq \delta(M) \ll_{\delta} M$, $\text{Rad}(H_i) \leq \delta(H_i) \neq H_i$, for all $i \in I$. Hence H_i is local for all $i \in I$.

(3) $\Rightarrow$ (2) Let A be a proper submodule of M . By Proposition 3.6, $A = N \oplus S$ where $N \leq^d M$ and $S \ll_{\delta} M$. Now M has a decomposition $M = \oplus_{j \in J} M_j$ with each M_j local. It then follows by Theorem 4.9, that this decomposition complements direct summands. Hence $M = N \oplus (\oplus_{k \in K} M_k)$, for some $K \subseteq J$. For a fixed $\alpha \in K$, it is clear that $B = N \oplus (\oplus_{k \neq \alpha} M_k) \oplus \delta(M_{\alpha})$ is a maximal submodule of M . Now $A = N \oplus S \leq N + \delta(M) \leq B$. Hence (2) follows.

(2) $\Rightarrow$ (1) Is trivial. □

Corollary 4.12 *Let M be an amply supplemented quasi δ -discrete module such that every δ -supplement submodule is δ -coclosed. Then M has a decomposition unique up to isomorphism, $M = M_1 \oplus M_2$, where $\delta(M_1) \ll_{\delta} M_1$ and $\delta(M_2) = M_2$.*

Proof By Theorem 4.9, $M = \oplus_{i \in I} H_i$ with each H_i hollow. Define $J = \{j \in I : \delta(H_j) \neq H_j\}$ and let $K = I \setminus J$. Let $M_1 = \oplus_{j \in J} H_j$ and $M_2 = \oplus_{k \in K} H_k$. Since H_j is local for all $j \in J$, by Proposition 4.11, $\delta(M_1) \ll_{\delta} M_1$. Now $\delta(M_2) = \oplus_{k \in K} \delta(H_k) = \oplus_{k \in K} H_k = M_2$. The uniqueness of the decomposition $M = M_1 \oplus M_2$ follows from the uniqueness of the original decomposition $M = \oplus_{i \in I} H_i$. □

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TOPOLOGICAL COLIMIT PRESERVING EQUALIZERS

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Abstract

For an exact category, $\mathcal{E}$, the colimit functor $Colim_C : \mathcal{E}^C \rightarrow \mathcal{E}$ preserves finite limits if and only if C is filtered. In previous work we investigated preservation of equalizers in case the base category is the non-exact category, Top , of topological spaces and showed that for a filtered internal category C in Top , topological colimit preserves equalizers if and only if C is admissible. In this manuscript we present conditions on an internal category C , making it admissible.

Keywords: colimit, equalizer, filtered internal category, topological spaces

AMS Subject Classification Codes: 18A30, 18A35, 18B30

1. Introduction and Preliminaries

For an exact category (hence for a topos) $\mathcal{E}$, the functor $Colim : \mathcal{E}^C \rightarrow \mathcal{E}$ is known to preserve finite limits if and only if C is a filtered internal category in $\mathcal{E}$, see [1], [4], and [7]. Others have worked on this in some special cases, see [3] and [5].

Since the category $KTop$ of compactly generated topological spaces is exact, see [8], the functor $Colim : KTop^C \rightarrow KTop$ preserves finite limits, if and only if C is an internal filtered category in $KTop$. The importance of this in algebraic topology stems from the fact that $Colim : KTop^{\Delta^*} \rightarrow KTop$ is part of the geometric realization functor $|-| : Set^{\Delta^{op}} \rightarrow KTop$. In fact the geometric realization functor $|-| : Set^{\Delta^{op}} \rightarrow KTop$ splits to three functors $Set^{\Delta^{op}} \xrightarrow{\bar{D}} KTop^{\Delta^{op}} \xrightarrow{O^*} KTop^{\Delta^*} \xrightarrow{Colim} KTop$, see [5], in which $\bar{D}$ and O^* are known to preserve finite limits and since the space Δ^* is filtered, $Colim$, preserves finite limits as well, and hence so does $|-|$.

The fact that the geometric realization functor preserves finite limits means that the pair $(S, |-|)$ is a geometric morphism, see [8] page 455, where $|-| : Set^{\Delta^{op}} \rightarrow KTop$ is the singular functor.

In [8], page 460, this situation has been generalized to the case where $KTop$ is replaced by a cocomplete topos $\mathcal{E}$. So the problem of when does the functor $Colim : A^C \rightarrow A$ preserve finite limits for a non-exact category A , generalizes the known theorems and yields a geometric morphism in a more generalized setting.

In [6], we investigated the problem of $Colim$ preserving equalizers, when the base category is the non-exact category, Top , of topological spaces. we presented a necessary and sufficient condition, on topological presheaves, for preservation of equalizers and defined the notion of an admissible category. In [9] we gave sufficient conditions for admissibility. In this article we present further conditions, on an internal category C , rendering admissibility.

Definition 1.1 ([6]). Let $X_1 \xrightarrow[d_1]{d_0} X_0 \xrightarrow{q} Y$ be maps in Set and H a subset of X_0 . Denoting the complement of a subset A of X_0 by A^c , we define subsets $\tilde{H}, \bar{H}$ of X_0 by $\tilde{H} = \bigcup_{A \in \mathcal{B}} A$, where $\mathcal{B} = \{A | A \subseteq H, d_0^{-1}(A) \subseteq d_1^{-1}(A)\}$, and $\bar{H} = (q^{-1}q(\tilde{H} \cup \tilde{H}^c))^c$.

Definition 1.2 ([6]). An internal category C in Top is said to be *ad-*

missible if for every discrete opfibration $X \xrightarrow{\gamma} C$ in Top^C and every open subset H of X_0 , there exists a subset $\hat{H}$ of $\tilde{H}$, such that $q^{-1}q(\tilde{H} \cup \hat{H})$ is an open subset of X_0 , where q is the coequalizer of the parallel pair $X_1 \xrightarrow[d_1]{d_0} X_0$.

Remark 1.3. The definition of filtered category given in [8] is in a topos. Since we prefer to start with the least condition on C , we use the same definition of filteredness in Top , rather than using the stronger conditions obtained by replacing "epimorphism" with "quotient map" or the like.

Theorem 1.4 ([6]). Let C be a filtered category in Top . Then the functor $Colim_C : Top^C \rightarrow Top$ preserves equalizers if and only if C is admissible.

2. Admissible category with finite topology

In this section, we impose conditions on an internal category with finite topology, (i.e., with finite topology on the object of morphisms), rendering it admissible.

Lemma 2.1 ([6]). Let X be an internal category in Set , q be a coequalizer of the parallel pair $X_1 \xrightarrow[d_1]{d_0} X_0$, and H be a subset of X_0 . Then

- (a) $\tilde{H} = \{x \in X_0 \mid d_0^{-1}(x) \subseteq d_1^{-1}(H)\}$.
- (b) $\tilde{H}^c = (d_0 d_1^{-1}(H))^c$.
- (c) If in addition, X is weakly filtered, then $q^{-1}q(\tilde{H}) = d_0 d_1^{-1}(\tilde{H})$. Also if $\tilde{H} = H$ or $\tilde{H}^c = H^c$, then $\tilde{H} = \emptyset$.

Lemma 2.2 ([6]). Let $X \xrightarrow{\gamma} C$ in $Cat(Set)$ be a discrete opfibration and D a subset of C_0 . Then

- (a) $\gamma_0^{-1}(D) = \gamma_0^{-1}(\tilde{D})$.
- (b) If $H = \gamma_0^{-1}(D)$, then $q_X^{-1}q_X(\tilde{H}) = \gamma_0^{-1}(q_C^{-1}q_C(\tilde{D}))$, where $q_X = coeq(X)$ and $q_C = coeq(C)$.
- (c) If C is filtered, $H = \gamma_0^{-1}(D)$ and $\tilde{D} \neq \emptyset$, then $q_X^{-1}q_X(\tilde{H}) = X_0$.

Remark 2.3. If C is an internal category in Set , then we have a discrete opfibration $X \xrightarrow{\gamma} C$, with $X_0 = C_1$, $X_1 = C_2$, $d_0 = \pi_1$, $d_1 = m$, $m = \pi_1 \pi_1 \times m(\pi_1 \pi_1 \times \pi_2 \pi_2)$. Also by Definition 1.1 and Lemma 2.1, for every subset B of C_1 we have:

$$\begin{aligned} \tilde{B} &= \{f \in B \mid \pi_1^{-1}(f) \subseteq m^{-1}(B)\} = \\ &= \{f \in B \mid gf \in B, \text{ for every } g \in C_1, \text{ composable with } f\} \end{aligned}$$

Definition 2.4. Let C be an internal category in Top . We say d_0 is $\sim$ -open, if for every open subset B of C_1 , $d_0(\tilde{B})$ is an open subset of C_0 and we say d_0 is $\sim$ -closed, if for every closed subset F of C_1 , $d_0(\tilde{F})$ is a closed subset of C_0 .

Definition 2.5. An internal category C in Top is said to be *induced*, if τ_{C_1} is the induced topology by $(d_0, d_1; \tau_{C_0})$.

Theorem 2.6 ([6]). Let C be a filtered category in Top , such that τ_{C_1} is a finite set.

- (a) If d_0 is $\sim$ -open, or $\sim$ -closed, then C is admissible.
- (b) If C is an induced filtered category, then C is admissible.

Definition 2.7. A map $A \xrightarrow{f} B$ in Top is said to be a *universal open (closed) map*, if the pullback of f along every map is an open (closed) map.

Lemma 2.8 A map $A \xrightarrow{f} B$ in Top is:

- (a) a universal open map if and only if it is an open map.
- (b) a universal closed map if and only if it is a proper map (i.e., the product map $Z \times X \xrightarrow{i_Z \times f} Z \times Y$ is closed for every space Z , where $Z \xrightarrow{i_Z} Z$ is the identity map) if and only if f is closed and the set $f^{-1}\{y\}$ is compact for every point $y \in Y$.

Lemma 2.9. Let C be a filtered category. If d_0 is an open map and a closed map, then d_0 is $\sim$ -closed.

Proof. Let F be a closed subset of C_1 . By Lemma 2.1, $\tilde{F} = (\pi_1 m^{-1}(F^c))^c$ and by Lemma 2.8, π_1 is an open map. So $d_0(\tilde{F})$ is a closed subset of C_0 $\square$

Corollary 2.10. Let C be a filtered category. If τ_{C_1} is finite, d_0 is an open map and a closed map, then C is admissible.

Proof. By Theorem 2.6 and Lemma 2.9, the result follows. $\square$

Example 2.11. Let $C_0 = \{a, b\}$ and $C_1 = \{1_a, 1_b, a \xrightarrow{\mu} b\}$ be the set of objects and the set of morphisms of the category C , respectively. Then C is filtered in Set . Let $\tau_{C_0} = \{\emptyset, C_0, \{a\}\}$ and $\tau_{C_1} = \{\emptyset, C_1, \{1_a, \mu\}, \{1_a\}, \{\mu\}\}$. Then C is filtered in Top , d_0 is an open map and a closed map. So by Corollary 2.10, C is admissible.

3. Admissible category with arbitrary topology

In this section, we impose conditions on an internal category with arbitrary topology, making it admissible.

Lemma 3.1. Let $X \xrightarrow{\gamma} C \in Cat(Set)$ and H be a subset of X_0 , such that $d_1^{-1}(H) = \bigcup_{\alpha \in I} d_0^{-1}(A_\alpha) \cap \gamma_1^{-1}(B_\alpha)$, where A_α and B_α are subsets of X_0 and C_1 , respectively, for every α in I . Then

(a) $\tilde{H} = \bigcup_{x \in H} ((\bigcap_{\alpha \in I_x} A_\alpha) \cap \gamma_0^{-1} i^{-1}(\widetilde{B_x}))$, where $I_x = \{\alpha \in I \mid x \in A_\alpha\}$ and $B_x = \bigcup_{\alpha \in I_x} B_\alpha$.

(b) If $d_0^{-1}(a) \subseteq \bigcup_{\alpha \in I} B_\alpha$ implies there exists a finite subset J of I , such that $d_0^{-1}(a) \subseteq \bigcup_{\alpha \in J} B_\alpha$, then $\tilde{H} = \bigcup_{J \in \Lambda} ((\bigcap_{\alpha \in J} A_\alpha) \cap \gamma_0^{-1} i^{-1}(\widetilde{B_J}))$, where $B_J = \bigcup_{\alpha \in J} B_\alpha$ and $\Lambda = \{J \mid J \text{ is a finite set of } I\}$.

Proof. (a) First for every subset B of C_1 , we have $d^{-1}(a) \subseteq B \Leftrightarrow 1_a \in \tilde{B} \Leftrightarrow a \in i^{-1}(\tilde{B})$, so $\{a \mid d_0^{-1}(a) \subseteq B\} = i^{-1}(\tilde{B})$.

Now let $x \in H$ and $H_x = (\bigcap_{\alpha \in I_x} A_\alpha) \cap \gamma_0^{-1} i^{-1}(\widetilde{B_x})$, then $d_o^{-1}(H_x) \subseteq d_1^{-1}(H)$, so by Lemma 2.1, $H_x \subseteq \tilde{H}$. Conversely let $x \in \tilde{H}$, by Lemma 2.1, $d_o^{-1}(x) \subseteq d_1^{-1}(H)$. Now we have $d_o^{-1}(x) \subseteq \bigcup_{\alpha \in I_x} \gamma_1^{-1}(B_\alpha) \Rightarrow \gamma_1 d_o^{-1}(x) \subseteq \bigcup_{\alpha \in I_x} B_\alpha \Rightarrow d_0^{-1} \gamma_0(x) \subseteq \bigcup_{\alpha \in I_x} B_\alpha \Rightarrow \gamma_0(x) \in i^{-1}(\widetilde{B_x}) \Rightarrow x \in \gamma_0^{-1} i^{-1}(\widetilde{B_x})$. Therefore $x \in H_x$.

(b) Let $x \in \tilde{H}$, by (a) $x \in (\bigcap_{\alpha \in I_x} A_\alpha) \cap \gamma_0^{-1} i^{-1}(\widetilde{B_x})$. So $d_0^{-1}(\gamma_0(x)) \subseteq \bigcup_{\alpha \in I_x} B_\alpha$. By hypothesis, there exists a finite subset J of I_x , such that $d_0^{-1}(\gamma_0(x)) \subseteq \bigcup_{\alpha \in J} B_\alpha$. It follows that $\gamma_0(x) \in i^{-1}(\widetilde{B_J})$ and therefore $x \in (\bigcap_{\alpha \in J} A_\alpha) \cap \gamma_0^{-1} i^{-1}(\widetilde{B_J})$.

Conversely, by the proof of (a) we have $(\bigcap_{\alpha \in J} A_\alpha) \cap \gamma_0^{-1} i^{-1}(\widetilde{B_J}) \subseteq \tilde{H}$. $\square$

Let X be an internal category. A subset H of X_0 is (d_0, d_1) -semisaturated if $d_o^{-1}(H) \subseteq d_1^{-1}(H)$. By Definition 1.1, for every subset H of X_0 , $\tilde{H}$ is (d_0, d_1) -semisaturated.

Definition 3.2. Let C be an internal category in Top and $\widetilde{\tau_{C_1}}$ be the

topology generated by the basis $\{\tilde{B} \mid B \in \tau_{C_1}\}$.

(a) We say C has *weak semisaturated topology*, if $\widetilde{\tau_{C_1}} \subseteq \tau_{C_1}$.

(b) We say C has *semisaturated topology*, if $\widetilde{\tau_{C_1}} = \tau_{C_1}$.

(c) We say C is *compactly sourced*, if for every $a \in C_0$, $d_0^{-1}(a)$ is a compact subset of C_1 with respect to τ_{C_1} .

Lemma 3.3. Let C be an internal category in Top with semisaturated topology. Then C has weak semisaturated topology and is compactly sourced.

Proof. It is obvious that C has weak semisaturated topology. Now let $a \in C_0$ and $d_0^{-1}(a) \subseteq \bigcup_{\alpha \in I} B_\alpha$, where B_α is an open subset of C_1 , for every $\alpha \in I$. Then $1_a \in B_\alpha$ for some $\alpha \in I$. Since C has semisaturated topology, so $1_a \in B_\alpha = \widetilde{B}_\alpha$. Therefore $d_0^{-1}(a) \subseteq B_\alpha$. $\square$

Theorem 3.4. If C is filtered, has weak semisaturated topology, is compactly sourced, and $C_1 \xrightarrow{d_0} C_0$ is an open map, then C is admissible.

Proof. Let $X \xrightarrow{\gamma} C \in Top^C$ and H be an open subset of X_0 . Since X_1 is a pullback in Top , so τ_{X_1} is the induced topology by $(d_0, \tau_{X_0}; \gamma_1, \tau_{C_1})$. Therefore there exist open subset A_α of X_0 and B_α of C_1 , for an index set I , such that $d_1^{-1}(H) = \bigcup_{\alpha \in I} d_0^{-1}(A_\alpha) \cap \gamma_1^{-1}(B_\alpha)$. By Lemma 3.1 (b), $\tilde{H}$ is

an open subset of X_0 , so by Lemma 2.1, $q^{-1}q(\tilde{H}) = d_0 d_1^{-1}(\tilde{H})$ is an open subset of X_0 . $\square$

Corollary 3.5. If C is filtered, has semisaturated topology, and $C_1 \xrightarrow{d_0} C_0$ is an open map, then C is admissible.

Proof. By Theorem 3.4 and Lemma 3.3, the result follows. $\square$

Theorem 3.6. If C is filtered, $C_1 \xrightarrow{d_0} C_0$ is an open and a universally closed map, then C is admissible.

Proof. Let $X \xrightarrow{\gamma} C \in Top^C$ and H be an open subset of X_0 . By Lemma 2.1, $\widetilde{H^c} = (d_0 d_1^{-1}(H))^c$ is a closed subset of X_0 . Since d_0 is a universal closed map, so $q^{-1}q(\tilde{H} \cup \tilde{H}) = (q^{-1}q(\widetilde{H^c}))^c = (d_0 d_1^{-1}(\widetilde{H^c}))^c$ is an open set. $\square$

Definition 3.7 An internal category C in Top is said to be *source-compact* if for every open collection $\{D_\alpha : \alpha \in I\}$ of C_0 such that $\bigcup_{\alpha \in I} D_\alpha \neq \emptyset$, there exists a finite subcollection $\{D_\alpha : \alpha \in J\}$ such that $\bigcup_{\alpha \in J} D_\alpha \neq \emptyset$.

Theorem 3.8. Let C be a filtered category. If C is induced and source-compact, then C is admissible.

Proof. Let $X \xrightarrow{\gamma} C \in Top^C$ and H be an open subset of X_0 . Since X_1 is a pullback in Top , so τ_{X_1} is the induced topology by $(d_0, \tau_{X_0}; \gamma_1, \tau_{C_1})$. Therefore there exist open subsets A_α of X_0 and D_α of C_0 , for an index set I , such that $d_1^{-1}(H) = \bigcup_{\alpha \in I} d_0^{-1}(A_\alpha) \cap \gamma_1^{-1}d_1^{-1}(D_\alpha)$.

Let $\Lambda = \{J \mid \bigcup_{\alpha \in J} D_\alpha \neq \emptyset \text{ and } J \text{ is a finite set of } I\}$. Then we show that $q^{-1}q(\tilde{H}) = \bigcup_{J \in \Lambda} (\bigcap_{\alpha \in J} A_\alpha)$. Let J be a finite set of I . By Lemma 3.1, and the fact that for any set A , $d_1d_0^{-1}(A)$ is (d_0, d_1) -semisaturated, we have

$$\begin{aligned} d_0^{-1}(\bigcap_{\alpha \in J} A_\alpha) \cap d_1^{-1}\gamma_0^{-1}(\bigcup_{\alpha \in J} D_\alpha) &\subseteq d_1^{-1}(H) \Rightarrow \\ d_1d_0^{-1}(\bigcap_{\alpha \in J} A_\alpha) \cap \gamma_0^{-1}(\bigcup_{\alpha \in J} D_\alpha) &\subseteq H \Rightarrow \\ d_1d_0^{-1}(\bigcap_{\alpha \in J} A_\alpha) \cap \gamma_0^{-1}(\bigcup_{\alpha \in J} \widetilde{D_\alpha}) &\subseteq \tilde{H} \Rightarrow \\ \bigcap_{\alpha \in J} (A_\alpha) &\subseteq q^{-1}q(\tilde{H}) \end{aligned}$$

Conversely let $x \in q^{-1}q(\tilde{H}) = d_0d_1^{-1}(\tilde{H})$, then there exists a morphism $x \xrightarrow{(x, \mu)} y$ in X_1 such that $y \in \tilde{H}$. So $(x, \mu) \in d_1^{-1}(H)$ and hence $\gamma_0(y) \in \bigcup_{\alpha \in I_x} D_\alpha$, where $I_x = \{\alpha \in I \mid x \in A_\alpha\}$. Then we show that $\bigcup_{\alpha \in I_x} D_\alpha \neq \emptyset$.

Let $\bigcup_{\alpha \in I_x} \widetilde{D_\alpha} = \emptyset$, so there exists a morphism $y \xrightarrow{(y, \mu')} z$ in X_1 such that $z \notin \gamma_0^{-1}(\bigcup_{\alpha \in I_x} D_\alpha)$. Since $\tilde{H}$ is a (d_0, d_1) -semisaturated subset of X_0 , so $z \in \tilde{H}$ and hence the morphism $(x, \mu'\mu) \in d_1^{-1}(H)$. Thus $z \in \gamma_0^{-1}(\bigcup_{\alpha \in I_x} D_\alpha)$, and

this is a contradiction. So $\bigcup_{\alpha \in I_x} \widetilde{D_\alpha} \neq \emptyset$, and by assumption there exists a finite set J of I_x , such that $\bigcup_{\alpha \in J} D_\alpha \neq \emptyset$. Hence $x \in \bigcap_{\alpha \in I_x} A_\alpha \subseteq \bigcap_{\alpha \in J} A_\alpha$. $\square$

Example 3.9. Let C be an internal category in Set , where the objects of C are natural numbers and there is a unique morphism $m \rightarrow n$ if and only

if $m \geq n$. Since C has terminal object 1, so it is a filtered category in *Set*. Let τ_{C_0} be the topology generated by the basis $\{\{1\}, \{1, 2\}, \{1, 2, 3\}, \dots\}$ and τ_{C_1} be the induced topology by (d_0, d_1, τ_{C_0}) . Then C is a filtered category in *Top*. Since for every open subset D of C_0 , $\tilde{D} = D$, C is source-compact. So by Theorem 3.8, C is admissible.

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**UNIFICATION OF q -FRACTIONAL
INTEGRAL OPERATORS**

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Abstract

In this paper a unification of Erdélyi-Kober and Weyl fractional q -integral operators is defined, further some rules of composition and the q -Mellin transform for these new operators are presented. Various particular cases are obtained and results given earlier by Galué can be derived from our results.

Key words: Erdélyi-Kober fractional q -integral operators, Weyl fractional q -integral operator, q -Mellin transform, composition rules.

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1 Introduction

The concept of differintegral of complex order ν , which is a generalization of the ordinary n -th derivative and n times integral to any complex number, can be introduced in several ways. The most widely used definition of an integral of fractional order is via an integral transform, called the Riemann-Liouville operator of fractional integration: [17, p. 146]

$$\begin{aligned} {}_a I_x^\alpha \varphi(x) &= \frac{1}{\Gamma(\alpha)} \int_a^x (x-t)^{\alpha-1} \varphi(t) dt \quad \operatorname{Re}(\alpha) > 0, \\ &= \frac{d^n}{dx^n} {}_a I_x^{\alpha+n} \varphi(x) \quad -n < \operatorname{Re}(\alpha) \leq 0, \quad n \in \mathbb{N}. \end{aligned} \quad (1)$$

Many authors ([1]-[3], [5]-[7], [9]-[14], [17]-[21]) have defined and studied operators of fractional integration through an integral transform. Some of these operators are:

1.1 Erdélyi-Kober Operator: [9, p. 4, No. (20)]

$$\begin{aligned} I_{\eta, \alpha} f(x) &= \frac{x^{-\alpha-\eta}}{\Gamma(\alpha)} \int_0^x (x-t)^{\alpha-1} t^\eta f(t) dt, \quad \operatorname{Re}(\alpha) \geq 0, \\ &= x^{-\alpha-\eta} \frac{d^n}{dx^n} x^{\eta+\alpha+n} I_{\eta, \alpha+n} f(x), \quad \operatorname{Re}(\alpha) < 0, \end{aligned} \quad (2)$$

where n is the minor integer major that α .

1.2 Basic analogue of Kober fractional q -integral operator

A basic analogue of Kober fractional q -integral operator has been defined by Agarwal [1] in the following form:

$$I_q^{\eta, \mu} f(x) = \frac{x^{-(\eta+\mu)}}{\Gamma_q(\mu)} \int_0^x t^\eta (x-tq)_{\mu-1} f(t) d_q t, \quad \eta \in \mathbb{C}, \operatorname{Re}(\mu) > 0, \quad (3)$$

where μ is an arbitrary order of integration, and

$$(x-y)_\nu = x^\nu \prod_{k=0}^{\infty} \frac{(1 - (\frac{y}{x}) q^k)}{(1 - (\frac{y}{x}) q^{k+\nu})}. \quad (4)$$

1.3 Basic analogue of Weyl fractional q -integral operator

Al-Salam [2] introduced a basic analogue of Weyl fractional q -integral operator as follows:

$$K_q^{\eta, \mu} f(x) = \frac{q^{-\eta} x^\eta}{\Gamma_q(\mu)} \int_x^\infty (t-x)_{\mu-1} t^{-(\eta+\mu)} f(q^{(1-\mu)} t) d_q t, \quad \eta \in \mathbb{C}, \operatorname{Re}(\mu) > 0. \quad (5)$$

1.4 The generalized Erdélyi-Kober fractional q -integral operator: Defined by [6]

$$\begin{aligned} I_q^{\eta, \mu, \beta} f(x) &= \frac{\beta x^{-\beta(\eta+\mu)}}{\Gamma_q(\mu)} \int_0^x (x^\beta - t^\beta q)_{\mu-1} t^{\beta(\eta+1)-1} f(t) d_q t \\ &= \beta(1 - q^{1/\beta})(1 - q)^{\mu-1} \sum_{k=0}^{\infty} \frac{(q^\mu; q)_k}{(q; q)_k} q^{k(\eta+1)} f(xq^{k/\beta}) \quad (6) \end{aligned}$$

$$\operatorname{Re}(\beta), \operatorname{Re}(\mu) > 0, \quad \eta \in \mathbb{C}.$$

1.5 The generalized Weyl fractional q -integral operator: [7]

$$\begin{aligned} K_q^{\eta, \mu, \beta} f(x) &= \frac{\beta q^{-\eta} x^{\beta\eta}}{\Gamma_q(\mu)} \int_x^\infty (t^\beta - x^\beta)_{\mu-1} t^{-\beta(\eta+\mu-1)-1} f(q^{(1-\mu)/\beta} t) d_q t \\ &= \beta(1 - q^{1/\beta})(1 - q)^{\mu-1} \sum_{k=0}^{\infty} \frac{(q^\mu; q)_k}{(q; q)_k} q^{k\eta} f(xq^{-(\mu+k)/\beta}) \quad (7) \end{aligned}$$

$$\operatorname{Re}(\beta), \operatorname{Re}(\mu) > 0, \quad \eta \in \mathbb{C}.$$

Operators of fractional calculus and their q -analogues have many applications, for example, they can be used to solve dual integral and series equations which arise in crack problems in elasticity [22]. One important property of fractional integrals is that certain multiples of them map some orthogonal polynomials to orthogonal polynomials [15].

In this paper a unification of Erdélyi-Kober and Weyl fractional q -integral operators is defined, further some rules of composition and

the q -Mellin transform for these new operators are presented. Various particular cases are obtained and results given earlier by Galué can be derived from our results.

2 Basic Hypergeometric Series

In this section we present some definitions necessary for the development of the next sections.

2.1 The q -shifted factorial

The q -shifted factorial is defined as: [8]

$$(a; q)_n = \begin{cases} 1, & n = 0 \\ (1-a)(1-aq)(1-aq^2)\dots(1-aq^{n-1}), & n = 1, 2, \dots \\ [(1-aq^{-1})(1-aq^{-2})\dots(1-aq^{-n})]^{-1}, & n = -1, -2, \dots \end{cases} \quad (8)$$

Also,

$$(a; q)_\infty = \lim_{n \rightarrow \infty} (a; q)_n = \prod_{k=0}^{\infty} (1 - aq^k), \quad (9)$$

which converges for $|q| < 1$ and diverges for $a \neq 0$ and $|q| \geq 1$, and

$$(a; q)_n = \frac{(a; q)_\infty}{(aq^n; q)_\infty}, \quad n \in \mathbb{Z}, |q| < 1. \quad (10)$$

2.2 Identities

We recall here the following q -identities given by Gasper and Rahman [8, p. 20, Nos. 1.2, 1.3(i); p. 235, No. (I.35)]:

$$(ab; q)_n = \sum_{k=0}^n \begin{bmatrix} n \\ k \end{bmatrix}_q (a; q)_k (b; q)_{n-k} b^k, \quad (11)$$

where

$$\begin{aligned} \begin{bmatrix} n \\ k \end{bmatrix}_q &= \frac{(q; q)_n}{(q; q)_k (q; q)_{n-k}}. \\ \Gamma_q(x) &= \frac{(q; q)_\infty}{(q^x; q)_\infty} (1-q)^{1-x}, \quad 0 < q < 1. \end{aligned} \quad (12)$$

From this result we can write

$$\frac{(q^y; q)_\infty}{(q^x; q)_\infty} = \frac{\Gamma_q(x)}{\Gamma_q(y)} (1 - q)^{x-y}. \quad (13)$$

2.3 The q -binomial theorem

One of the most important summation formulas for hypergeometric series is given by the binomial theorem:

$${}_2F_1(a, c; c; z) = {}_1F_0(a; -; z) = \sum_{n=0}^{\infty} \frac{(a)_n}{n!} z^n = (1 - z)^{-a}, \quad |z| < 1, \quad (14)$$

whose q -analogue, derived by Cauchy (1843), Heine (1847) and other mathematicians, is [8]

$${}_1\phi_0(a; -; q, z) = \sum_{n=0}^{\infty} \frac{(a; q)_n}{(q; q)_n} z^n = \frac{(az; q)_\infty}{(z; q)_\infty}, \quad |z| < 1, |q| < 1. \quad (15)$$

2.4 The q -derivative operator

It is denoted by $\mathcal{D}_q$ and defined for fixed q by [8, p. 22]

$$\mathcal{D}_q f(z) = \frac{f(z) - f(qz)}{(1 - q)z}. \quad (16)$$

2.5 The q -integral

Thomae (1869) and Jackson (1910) introduced the q -integral by [8, p. 19, No. (1.11.1)]

$$\int_0^1 f(t) d_q t = (1 - q) \sum_{n=0}^{\infty} f(q^n) q^n. \quad (17)$$

Jackson also defined an integral on $(0, \infty)$ by

$$\int_0^\infty f(t) d_q t = (1 - q) \sum_{n=-\infty}^{\infty} f(q^n) q^n, \quad (18)$$

then from (17) and (18) we have

$$\int_1^\infty f(t) d_q t = (1 - q) \sum_{n=1}^{\infty} f(q^{-n}) q^{-n}. \quad (19)$$

3 Definition of the Operator

In this section we introduce a generalized fractional q -integral operator, denoted by ${}_l\mathcal{J}_q^{\eta,\mu,\beta,\sigma}f(x)$, and defined as follows:

$${}_l\mathcal{J}_q^{\eta,\mu,\beta,\sigma}f(x) = \beta(1 - q^{1/\beta})(1 - q)^{\mu-1} \sum_{k=0}^{\infty} \frac{(q^\mu; q)_k}{(q; q)_k} q^{k(\eta+\ell)} f\left(xq^{(\sigma+(-1)^{\ell+1}k)/\beta}\right), \quad (20)$$

$\ell = 0, 1, 2, \dots$, $\operatorname{Re}(\beta), \operatorname{Re}(\mu) > 0$, $\eta, \sigma \in \mathbb{C}$, μ is an arbitrary order of integration.

As particular cases of this operator we have

$${}_{2\ell+1}\mathcal{J}_q^{\eta,\mu,\beta,\sigma}f(x) = I_q^{\eta+2\ell,\mu,\beta,\sigma}f(x) \quad (21)$$

where

$$I_q^{\eta,\mu,\beta,\sigma}f(x) = \frac{\beta x^{-\beta(\eta+\mu)}}{\Gamma_q(\mu)} \int_0^x (x^\beta - t^\beta q)_{\mu-1} t^{\beta(\eta+1)-1} f(q^{\sigma/\beta}t) d_q t \quad (22)$$

$$\operatorname{Re}(\beta), \operatorname{Re}(\mu) > 0, \eta, \sigma \in \mathbb{C},$$

and

$${}_{2\ell}\mathcal{J}_q^{\eta,\mu,\beta,\sigma}f(x) = K_q^{\eta+2\ell,\mu,\beta,\sigma}f(x) \quad (23)$$

where

$$K_q^{\eta,\mu,\beta,\sigma}f(x) = \frac{\beta q^{-\eta} x^{\beta\eta}}{\Gamma_q(\mu)} \int_x^\infty (t^\beta - x^\beta)_{\mu-1} t^{-\beta(\eta+\mu-1)-1} f(q^{(1+\sigma)/\beta}t) d_q t \quad (24)$$

$$\operatorname{Re}(\beta), \operatorname{Re}(\mu) > 0, \eta, \sigma \in \mathbb{C}.$$

Proof of (21) Applying the definition (22)

$$I_q^{\eta+2\ell,\mu,\beta,\sigma}f(x) = \frac{\beta x^{-\beta(\eta+2\ell+\mu)}}{\Gamma_q(\mu)} \int_0^x (x^\beta - t^\beta q)_{\mu-1} t^{\beta(\eta+2\ell+1)-1} f(q^{\sigma/\beta}t) d_q t,$$

making a change of variable and using (16), we get

$$I_q^{\eta+2\ell,\mu,\beta,\sigma} f(x) = \frac{\beta(1-q^{1/\beta})}{(1-q)} \frac{x^{-\beta(\mu-1)}}{\Gamma_q(\mu)} \int_0^1 (x^\beta - x^\beta yq)_{\mu-1} y^{\eta+2\ell} f(xq^{\sigma/\beta} y^{1/\beta}) d_q y. \quad (25)$$

Since from (4) and (9)

$$(x^\beta - x^\beta yq)_{\mu-1} = \frac{(yq; q)_\infty}{(yq^\mu; q)_\infty} x^{\beta(\mu-1)}$$

the result (25) is equivalent to

$$I_q^{\eta+2\ell,\mu,\beta,\sigma} f(x) = \frac{\beta(1-q^{1/\beta})}{(1-q)\Gamma_q(\mu)} \int_0^1 \frac{(yq; q)_\infty}{(yq^\mu; q)_\infty} y^{\eta+2\ell} f(xq^{\sigma/\beta} y^{1/\beta}) d_q y.$$

Now, in view of definition of q-integral given in (17) we can write

$$I_q^{\eta+2\ell,\mu,\beta,\sigma} f(x) = \frac{\beta(1-q^{1/\beta})}{\Gamma_q(\mu)} \sum_{k=0}^{\infty} \frac{(q^k q; q)_\infty}{(q^k q^\mu; q)_\infty} q^{k(\eta+2\ell+1)} f(xq^{(\sigma+k)/\beta}),$$

which by application of (10) and (12) leads us to the expression

$$I_q^{\eta+2\ell,\mu,\beta,\sigma} f(x) = \beta(1-q^{1/\beta})(1-q)^{\mu-1} \sum_{k=0}^{\infty} \frac{(q^\mu; q)_k}{(q; q)_k} q^{k(\eta+2\ell+1)} f(xq^{(\sigma+k)/\beta}).$$

The Interpretation of this result according to (20) complete the proof.

The result (23) is established in similar form as the former but using (19).

Other particular cases are

$${}_1\mathcal{J}_q^{\eta,\mu,\beta,0} f(x) = I_q^{\eta,\mu,\beta} f(x) \quad (26)$$

$${}_0\mathcal{J}_q^{\eta,\mu,\beta,-\mu} f(x) = K_q^{\eta,\mu,\beta} f(x) \quad (27)$$

where $I_q^{\eta,\mu,\beta} f(x)$ is the generalized Erdélyi-Kober fractional q -integral operator and $K_q^{\eta,\mu,\beta} f(x)$ is the generalized Weyl fractional q -integral operator given in (6) and (7) respectively.

Further

$$\lim_{q \rightarrow 1^-} {}_1\mathcal{J}_q^{\eta,\mu,\beta,0} f(x) = \lim_{q \rightarrow 1^-} I_q^{\eta,\mu,\beta} f(x) = I_\beta^{\eta,\mu} f(x) \quad (28)$$

$$\lim_{q \rightarrow 1^-} {}_0\mathcal{J}_q^{\eta,\mu,\beta,-\mu} f(x) = \lim_{q \rightarrow 1^-} K_q^{\eta,\mu,\beta} f(x) = K_\beta^{\eta,\mu} f(x) \quad (29)$$

where $\text{Re}(\beta), \text{Re}(\mu) > 0$, $\eta \in \mathbb{C}$, $I_\beta^{\eta,\mu}(\cdot)$ and $K_\beta^{\eta,\mu}$ are operators of Erdélyi-Kober [12, p. 147].

4 Composition formula for the ${}_e\mathcal{J}_q^{\eta,\mu,\beta,\sigma}(\cdot)$ fractional q -integral operator

In this section, we investigate various rules of composition for the ${}_e\mathcal{J}_q^{\eta,\mu,\beta,\sigma}(\cdot)$ fractional q -integral operator.

4.1 Composition formulas for two ${}_e\mathcal{J}_q^{\eta,\mu,\beta,\sigma}(\cdot)$ fractional q -integral operators

Theorem 1. *Let be ${}_e\mathcal{J}_q^{\eta,\mu,\beta,\sigma} f(x)$ a fractional q -integral operator, as defined in (20), then the following composition formulas hold*

$${}_e\mathcal{J}_q^{\eta_1,\mu_1,\beta,\sigma_1} {}_e\mathcal{J}_q^{\eta_2+\mu_1,\mu_2,\beta,\sigma_2} f(x) = \frac{\beta(1-q^{1/\beta})}{(1-q)} {}_e\mathcal{J}_q^{\eta_1,\mu_1+\mu_2,\beta,\sigma_1+\sigma_2} f(x). \quad (30)$$

$${}_e\mathcal{J}_q^{\eta_2-\mu_1,\mu_1,\beta,\sigma_1} {}_e\mathcal{J}_q^{\eta_2,\mu_2,\beta,\sigma_2} f(x) = \frac{\beta(1-q^{1/\beta})}{(1-q)} {}_e\mathcal{J}_q^{\eta_2-\mu_1,\mu_1+\mu_2,\beta,\sigma_1+\sigma_2} f(x). \quad (31)$$

$${}_e\mathcal{J}_q^{\eta_2+\mu_2,\mu_1,\beta,\sigma_1} {}_e\mathcal{J}_q^{\eta_2,\mu_2,\beta,\sigma_2} f(x) = \frac{\beta(1-q^{1/\beta})}{(1-q)} {}_e\mathcal{J}_q^{\eta_2,\mu_1+\mu_2,\beta,\sigma_1+\sigma_2} f(x). \quad (32)$$

$${}_e\mathcal{J}_q^{\eta_1,\mu_1,\beta,\sigma_1} {}_e\mathcal{J}_q^{\eta_1-\mu_2,\mu_2,\beta,\sigma_2} f(x) = \frac{\beta(1-q^{1/\beta})}{(1-q)} {}_e\mathcal{J}_q^{\eta_1-\mu_2,\mu_1+\mu_2,\beta,\sigma_1+\sigma_2} f(x), \quad (33)$$

where $\ell = 0, 1, 2, \dots, \operatorname{Re}(\beta), \operatorname{Re}(\mu_i) > 0, \eta_i, \sigma_i \in \mathbb{C}, i = 1, 2$.

Proof To prove the result (30), we make use of the relation (20) in the left hand side to obtain

$${}_ell\mathcal{J}_q^{\eta_1, \mu_1, \beta, \sigma_1} {}_ell\mathcal{J}_q^{\eta_1 + \mu_1, \mu_2, \beta, \sigma_2} f(x) = \beta^2 (1 - q^{1/\beta})^2 (1-q)^{\mu_1-1} (1-q)^{\mu_2-1} \times \\ \sum_{k_1=0}^{\infty} \frac{(q^{\mu_1}; q)_{k_1}}{(q; q)_{k_1}} q^{k_1(\eta_1 + \ell)} \sum_{k_2=0}^{\infty} \frac{(q^{\mu_2}; q)_{k_2}}{(q; q)_{k_2}} q^{k_2(\eta_1 + \mu_1 + \ell)} f \left(x q^{(\sigma_1 + \sigma_2 + (-1)^{\ell+1}(k_1 + k_2))/\beta} \right),$$

using the result [16, p. 56, No. 1]

$$\sum_{n=0}^{\infty} \sum_{k=0}^{\infty} A(k, n) = \sum_{n=0}^{\infty} \sum_{k=0}^n A(k, n - k)$$

we get

$${}_ell\mathcal{J}_q^{\eta_1, \mu_1, \beta, \sigma_1} {}_ell\mathcal{J}_q^{\eta_1 + \mu_1, \mu_2, \beta, \sigma_2} f(x) = \beta^2 (1 - q^{1/\beta})^2 (1-q)^{\mu_1-1} (1-q)^{\mu_2-1} \times \\ \sum_{k_1=0}^{\infty} \sum_{k_2=0}^{k_1} \frac{(q^{\mu_1}; q)_{k_1 - k_2}}{(q; q)_{k_1 - k_2}} \frac{(q^{\mu_2}; q)_{k_2}}{(q; q)_{k_2}} q^{k_1(\eta_1 + \ell)} q^{k_2 \mu_1} f \left(x q^{(\sigma_1 + \sigma_2 + (-1)^{\ell+1} k_1)/\beta} \right),$$

that is,

$${}_ell\mathcal{J}_q^{\eta_1, \mu_1, \beta, \sigma_1} {}_ell\mathcal{J}_q^{\eta_1 + \mu_1, \mu_2, \beta, \sigma_2} f(x) = \beta^2 (1 - q^{1/\beta})^2 (1-q)^{\mu_1-1} (1-q)^{\mu_2-1} \times \\ \sum_{k_1=0}^{\infty} q^{k_1(\eta_1 + \ell)} f \left(x q^{(\sigma_1 + \sigma_2 + (-1)^{\ell+1} k_1)/\beta} \right) \sum_{k_2=0}^{k_1} \frac{(q^{\mu_1}; q)_{k_1 - k_2}}{(q; q)_{k_1 - k_2}} \frac{(q^{\mu_2}; q)_{k_2}}{(q; q)_{k_2}} q^{k_2 \mu_1}$$

which by virtue of the formula (11) is equivalent to

$${}_ell\mathcal{J}_q^{\eta_1, \mu_1, \beta, \sigma_1} {}_ell\mathcal{J}_q^{\eta_1 + \mu_1, \mu_2, \beta, \sigma_2} f(x) = \beta^2 (1 - q^{1/\beta})^2 (1-q)^{\mu_1-1} (1-q)^{\mu_2-1} \times \\ \sum_{k_1=0}^{\infty} \frac{(q^{\mu_1 + \mu_2}; q)_{k_1}}{(q; q)_{k_1}} q^{k_1(\eta_1 + \ell)} f \left(x q^{(\sigma_1 + \sigma_2 + (-1)^{\ell+1} k_1)/\beta} \right).$$

Interpreting the above expression with the help of the relation (20), we finally arrive to the result.

The proof of the other rules of composition is analogous.

4.2 Composition formulas for n ${}_q J_q^{\eta, \mu, \beta, \sigma}(\cdot)$ fractional q -integral operators

The composition rules earlier can be extended for n operators in the following form:

Theorem 2 *Let be ${}_q J_q^{\eta, \mu, \beta, \sigma} f(x)$ a fractional q -integral operator, as defined in (20), then the following composition formulas hold*

$$\begin{aligned} & {}_q J_q^{\eta_1, \mu_1, \beta, \sigma_1} {}_q J_q^{\eta_1 + \mu_1, \mu_2, \beta, \sigma_2} \dots {}_q J_q^{\eta_1 + \mu_1 + \mu_2 + \dots + \mu_{n-2}, \mu_{n-1}, \beta, \sigma_{n-1}} \\ & {}_q J_q^{\eta_1 + \mu_1 + \mu_2 + \dots + \mu_{n-1}, \mu_n, \beta, \sigma_n} f(x) = \\ & \left[\frac{\beta (1 - q^{1/\beta})}{(1 - q)} \right]^{n-1} {}_q J_q^{\eta_1, \mu_1 + \mu_2 + \dots + \mu_n, \beta, \sigma_1 + \sigma_2 + \dots + \sigma_n} f(x). \end{aligned} \quad (34)$$

$$\begin{aligned} & {}_q J_q^{\eta_n - \mu_{n-1} - \dots - \mu_2 - \mu_1, \mu_1, \beta, \sigma_1} {}_q J_q^{\eta_n - \mu_{n-1} - \dots - \mu_2, \mu_2, \beta, \sigma_2} \dots {}_q J_q^{\eta_n - \mu_{n-1}, \mu_{n-1}, \beta, \sigma_{n-1}} \\ & {}_q J_q^{\eta_n, \mu_n, \beta, \sigma_n} f(x) = \\ & \left[\frac{\beta (1 - q^{1/\beta})}{(1 - q)} \right]^{n-1} {}_q J_q^{\eta_n - \mu_{n-1} - \dots - \mu_2 - \mu_1, \mu_1 + \mu_2 + \dots + \mu_n, \beta, \sigma_1 + \sigma_2 + \dots + \sigma_n} f(x). \end{aligned} \quad (35)$$

$$\begin{aligned} & {}_q J_q^{\eta_n + \mu_n + \dots + \mu_2, \mu_1, \beta, \sigma_1} {}_q J_q^{\eta_n + \mu_n + \dots + \mu_3, \mu_2, \beta, \sigma_2} \dots {}_q J_q^{\eta_n + \mu_n, \mu_{n-1}, \beta, \sigma_{n-1}} \\ & {}_q J_q^{\eta_n, \mu_n, \beta, \sigma_n} f(x) = \\ & \left[\frac{\beta (1 - q^{1/\beta})}{(1 - q)} \right]^{n-1} {}_q J_q^{\eta_n, \mu_1 + \mu_2 + \dots + \mu_n, \beta, \sigma_1 + \sigma_2 + \dots + \sigma_n} f(x). \end{aligned} \quad (36)$$

$${}_q J_q^{\eta_1, \mu_1, \beta, \sigma_1} {}_q J_q^{\eta_1 - \mu_2, \mu_2, \beta, \sigma_2} \dots {}_q J_q^{\eta_1 - \mu_{n-1} - \mu_{n-2} - \dots - \mu_2, \mu_{n-1}, \beta, \sigma_{n-1}}$$

$$\begin{aligned} & {}_{\ell}\mathcal{J}_q^{\eta_1-\mu_n-\mu_{n-1}-\dots-\mu_2,\mu_n,\beta,\sigma_n} f(x) = \\ & \left[\frac{\beta(1-q^{1/\beta})}{(1-q)} \right]^{n-1} {}_{\ell}\mathcal{J}_q^{\eta_1-\mu_n-\mu_{n-1}-\dots-\mu_2,\mu_1+\mu_2+\dots+\mu_n,\beta,\sigma_1+\sigma_2+\dots+\sigma_n} f(x), \end{aligned} \quad (37)$$

where $\ell = 0, 1, 2, \dots, \operatorname{Re}(\beta), \operatorname{Re}(\mu_i) > 0, \eta_i, \sigma_i \in \mathbb{C}, i = 1, \dots, n$.

Proof of (34): We employ mathematical induction principle with an observation that for $n = 2$ in (34), it yields to equation (30) as:

$${}_{\ell}\mathcal{J}_q^{\eta_1,\mu_1,\beta,\sigma_1} {}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1,\mu_2,\beta,\sigma_2} f(x) = \frac{\beta(1-q^{1/\beta})}{(1-q)} {}_{\ell}\mathcal{J}_q^{\eta_1,\mu_1+\mu_2,\beta,\sigma_1+\sigma_2} f(x). \quad (38)$$

We suppose that (34) is valid for $n = k - 1$, that is,

$$\begin{aligned} & {}_{\ell}\mathcal{J}_q^{\eta_1,\mu_1,\beta,\sigma_1} {}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1,\mu_2,\beta,\sigma_2} \dots {}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1+\mu_2+\dots+\mu_{k-3},\mu_{k-2},\beta,\sigma_{k-2}} \\ & {}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1+\mu_2+\dots+\mu_{k-2},\mu_{k-1},\beta,\sigma_{k-1}} f(x) = \\ & \left[\frac{\beta(1-q^{1/\beta})}{(1-q)} \right]^{k-2} {}_{\ell}\mathcal{J}_q^{\eta_1,\mu_1+\mu_2+\dots+\mu_{k-1},\beta,\sigma_1+\sigma_2+\dots+\sigma_{k-1}} f(x). \end{aligned} \quad (39)$$

On operating both sides of the relation (39) by the operator ${}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1+\mu_2+\dots+\mu_{k-1},\mu_k,\beta,\sigma_k}(\cdot)$, we obtain

$$\begin{aligned} & {}_{\ell}\mathcal{J}_q^{\eta_1,\mu_1,\beta,\sigma_1} {}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1,\mu_2,\beta,\sigma_2} \dots {}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1+\mu_2+\dots+\mu_{k-2},\mu_{k-1},\beta,\sigma_{k-1}} \\ & {}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1+\mu_2+\dots+\mu_{k-1},\mu_k,\beta,\sigma_k} f(x) = \left[\frac{\beta(1-q^{1/\beta})}{(1-q)} \right]^{k-2} \times \\ & {}_{\ell}\mathcal{J}_q^{\eta_1,\mu_1+\mu_2+\dots+\mu_{k-1},\beta,\sigma_1+\sigma_2+\dots+\sigma_{k-1}} {}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1+\mu_2+\dots+\mu_{k-1},\mu_k,\beta,\sigma_k} f(x). \end{aligned}$$

On applying the relation (38) to the right hand side of the above expression, we finally obtain:

$${}_{\ell}\mathcal{J}_q^{\eta_1,\mu_1,\beta,\sigma_1} {}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1,\mu_2,\beta,\sigma_2} \dots {}_{\ell}\mathcal{J}_q^{\eta_1+\mu_1+\mu_2+\dots+\mu_{k-2},\mu_{k-1},\beta,\sigma_{k-1}}$$

$${}_l\mathcal{J}_q^{\eta_1+\mu_1+\mu_2+\dots+\mu_{k-1},\mu_k,\beta,\sigma_k}f(x)=\left[\frac{\beta(1-q^{1/\beta})}{(1-q)}\right]^{k-1} \times$$

$${}_l\mathcal{J}_q^{\eta_1,\mu_1+\mu_2+\dots+\mu_{k-1}+\mu_k,\beta,\sigma_1+\sigma_2+\dots+\sigma_{k-1}+\sigma_k}f(x)$$

which is valid for $n = k$ ($k = 1, 2, \dots, n$). This completes the proof of (34).

The results (35)-(37) can be proved in similar manner.

4.3 Particular cases

It is clear, from the results (21), (23), (26)-(29) that theorem 1 and theorem 2 provide rules of composition for $I_q^{\eta+2\ell,\mu,\beta,\sigma}(\cdot)$, $K_q^{\eta+2\ell,\mu,\beta,\sigma}(\cdot)$, $I_q^{\eta,\mu,\beta}(\cdot)$, $K_q^{\eta,\mu,\beta}(\cdot)$, $I_\beta^{\eta,\mu}(\cdot)$, $K_\beta^{\eta,\mu}(\cdot)$ operators. See, for example, rules of composition for $I_q^{\eta,\mu,\beta}(\cdot)$ given in [6].

5 q-Mellin Transform

In this section we establish the q-Mellin transform of the ${}_l\mathcal{J}_q^{\eta,\mu,\beta,\sigma}(\cdot)$ q-fractional integral operator.

Theorem 3 *Let be ${}_l\mathcal{J}_q^{\eta,\mu,\beta,\sigma}f(x)$ a fractional q-integral operator, as defined in (20), then*

$$M_q\{x^\rho {}_l\mathcal{J}_q^{\eta,\mu,\beta,\sigma}f(x); s\} = \frac{\beta(1-q^{1/\beta})}{(1-q)}q^{-\sigma(\rho+s)/\beta} \times$$

$$\frac{\Gamma_q(\eta + \ell - (-1)^{\ell+1}(\rho + s)/\beta)}{\Gamma_q(\mu + \eta + \ell - (-1)^{\ell+1}(\rho + s)/\beta)} M_q\{f(x); \rho + s\}. \quad (40)$$

where $M_q\{f(x); s\}$ denote the q-Mellin transform of $f(x)$ and $\ell = 0, 1, 2, \dots, \text{Re}(\beta), \text{Re}(\mu) > 0, \eta, \sigma \in \mathbb{C}$.

Proof The q-Mellin Transform of $f(x)$ is defined by [4]

$$M_q\{f(x); s\} = \frac{1}{(1-q)} \int_0^\infty x^{s-1} f(x) d_q x, \quad (41)$$

therefore, from (20) and (41) we obtain,

$$M_q \{ x^\rho {}_\ell \mathcal{J}_q^{\eta, \mu, \beta, \sigma} f(x); s \} = \frac{\beta(1 - q^{1/\beta})(1 - q)^{\mu-1}}{(1 - q)} \times \sum_{k=0}^{\infty} \frac{(q^\mu; q)_k}{(q; q)_k} q^{k(\eta+\ell)} \int_0^\infty x^{\rho+s-1} f(xq^{(\sigma+(-1)^{\ell+1}k)/\beta}) d_q x, \quad (42)$$

where we have interchanged the order of integral and summation.

Making a change of variable and using (16) we have, after some simplifications,

$$\int_0^\infty x^{\rho+s-1} f(xq^{(\sigma+(-1)^{\ell+1}k)/\beta}) d_q x = q^{-\sigma(\rho+s)/\beta} q^{-(-1)^{\ell+1}(\rho+s)k/\beta} \times (1 - q) \int_0^\infty w^{\rho+s-1} f(w) d_q w,$$

that is,

$$\int_0^\infty x^{\rho+s-1} f(xq^{(\sigma+(-1)^{\ell+1}k)/\beta}) d_q x = q^{-\sigma(\rho+s)/\beta} q^{-(-1)^{\ell+1}(\rho+s)k/\beta} \times (1 - q) M_q \{ f(w); \rho + s \}.$$

The substitution of this expression in (42) leads us to

$$M_q \{ x^\rho {}_\ell \mathcal{J}_q^{\eta, \mu, \beta, \sigma} f(x); s \} = \beta(1 - q^{1/\beta})(1 - q)^{\mu-1} q^{-\sigma(\rho+s)/\beta} \times M_q \{ f(w); \rho + s \} \sum_{k=0}^{\infty} \frac{(q^\mu; q)_k}{(q; q)_k} q^{k(\eta+\ell-(-1)^{\ell+1}(\rho+s)/\beta)}. \quad (43)$$

On the other hand, applying (15) and (13)

$$\sum_{k=0}^{\infty} \frac{(q^\mu; q)_k}{(q; q)_k} q^{k(\eta+\ell-(-1)^{\ell+1}(\rho+s)/\beta)} = \frac{(q^{\mu+\eta+\ell-(-1)^{\ell+1}(\rho+s)/\beta}; q)_\infty}{(q^{\eta+\ell-(-1)^{\ell+1}(\rho+s)/\beta}; q)_\infty}$$

$$= \frac{\Gamma_q(\eta + \ell - (-1)^{\ell+1}(\rho + s)/\beta)}{\Gamma_q(\mu + \eta + \ell - (-1)^{\ell+1}(\rho + s)/\beta)} (1 - q)^{-\mu}.$$

From this result and (43) we get (40).

Particular cases:

i) From (26) and (40):

$$M_q \{x^\rho I_q^{\eta, \mu, \beta} f(x); s\} = \frac{\beta(1 - q^{1/\beta})}{(1 - q)} \frac{\Gamma_q(\eta + 1 - (\rho + s)/\beta)}{\Gamma_q(\mu + \eta + 1 - (\rho + s)/\beta)} \times \\ M_q \{f(x); \rho + s\}. \quad (44)$$

ii) From (27) and (40):

$$M_q \{x^\rho K_q^{\eta, \mu, \beta} f(x); s\} = \frac{\beta(1 - q^{1/\beta})}{(1 - q)} q^{\mu(\rho+s)/\beta} \frac{\Gamma_q(\eta + (\rho + s)/\beta)}{\Gamma_q(\mu + \eta + (\rho + s)/\beta)} \times \\ M_q \{f(x); \rho + s\}. \quad (45)$$

iii) If $\beta = 1$, $\rho = 0$ in (44) and (45) we obtain respectively:

$$M_q \{I_q^{\eta, \mu} f(x); s\} = \frac{\Gamma_q(\eta + 1 - s)}{\Gamma_q(\mu + \eta + 1 - s)} M_q \{f(x); s\}. \quad (46)$$

$$M_q \{K_q^{\eta, \mu} f(x); s\} = \frac{\Gamma_q(\eta + s) q^{\mu s}}{\Gamma_q(\mu + \eta + s)} M_q \{f(x); s\}. \quad (47)$$

iv) From (46) with $\eta = 0$:

$$M_q \{I_q^{0, \mu} f(x); s\} = M_q \{x^{-\mu} I_q^\mu f(x); s\} = \frac{\Gamma_q(1 - s)}{\Gamma_q(\mu + 1 - s)} M_q \{f(x); s\} \quad (48)$$

where $I_q^\mu(\cdot)$ is a q-analogue of the Riemann-Liouville fractional integral operator due to Agarwal [1].

v) If $\eta = -\mu$ in (47):

$$M_q \{K_q^{-\mu, \mu} f(x); s\} = M_q \{q^{\mu(\mu+1)/2} x^{-\mu} K_q^\mu f(x); s\} = \frac{\Gamma_q(s - \mu) q^{\mu s}}{\Gamma_q(s)} \times$$

$$M_q\{f(x); s\} \quad (49)$$

where $K_q^\mu(\cdot)$ is a q -analogue of the Weyl fractional integral operator due to Al-Salam [2].

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**STUDIES ON THE LIE-SANTILLI ISOTHEORY
WITH UNIT OF GENERAL FORM¹**

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Abstract

We review matrix Lie-Santilli algebras and groups with associative and distributive matrix product, i.e. the theory without assuming that the unit matrix has the conventional form. The axiom of associative and distributive matrix product can be realized in different ways implying accordingly different forms of the unit matrix (Santilli isounit) obeying the axiom of unit element. Such algebras were first studied by Santilli. We study realizations of Lie-Santilli groups and algebras, and review examples of Lie-Santilli matrix algebras which are of interest in physics.

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Contents

1	Introduction	2
2	Matrix algebra $M(n, \mathbb{C}, \hat{\times})$	5
2.1	The $\hat{\times}$ -product of matrices	5
2.2	The dual algebra $M(n, \mathbb{C}, \hat{\times}^{-1})$	8
2.3	Metrics and coordinates of $M(n, \mathbb{C}, \hat{\times})$	9
2.4	Homotopy class of unit	11
2.4.1	Path from $\hat{I}$ to I	11
2.4.2	$\hat{I}$ in $M(2, \mathbb{C})$	12
2.4.3	The generating matrix	13
3	Lie-Santilli groups and Lie-Santilli algebras	14
3.1	Lie-Santilli groups	14
3.1.1	Group $GL(n, \mathbb{C}, \hat{\times})$	14
3.1.2	Unitary group $U(n, \mathbb{C}, \hat{\times})$	14
3.1.3	Orthogonal group $O(n, \mathbb{R}, \hat{\times})$	15
3.1.4	Group $SL(n, \mathbb{C}, \hat{\times})$	15
3.2	Action of the groups on classical linear spaces	15
3.2.1	Eigenvalue problem	15
3.2.2	Action of unitary group	16
3.2.3	Action of orthogonal group	17
3.2.4	Action of pseudo-unitary and pseudo-orthogonal groups	19
3.3	Matrix exponent and Lie-Santilli algebras	20
3.4	Some restrictions on the form of Santilli isounit	24
3.5	Infinite-dimensional case	24
4	Examples of matrix Lie-Santilli algebras and groups	26
4.1	Algebra $so(3, \mathbb{R}, \hat{\times})$ and group $SO(3, \mathbb{R}, \hat{\times})$	26
4.2	The groups $SO(2, \mathbb{R}, \hat{\times})$ and $U(1, \mathbb{C}, \hat{\times})$	28
4.3	Action of the group $U(1, \mathbb{C}, \hat{\times})$	30
4.4	Group $SO(1, 1, \mathbb{R}, \hat{\times})$	33
4.5	Realization map	35
4.6	Matrix algebra $M(2, \mathbb{C}, \hat{\times})$	37
4.7	Algebra $su(2, \mathbb{C}, \hat{\times})$ and group $SU(2, \mathbb{C}, \hat{\times})$	39
A	Appendix A	41
B	Appendix B	43

1 Introduction

The unit matrix in $n \times n$ matrix algebra over field of complex numbers $\mathbb{C}$ is defined due to the axiom of unit element. This axiom uses matrix product which in turn is axiomatically defined as an associative and distributive one [1, 2]. The unit is taken as a diagonal $n \times n$ matrix of the following conventional form:

$$I = \text{diag}(1, 1, \dots, 1).$$

The matrix I is such that for any $n \times n$ matrix M the following relations hold:

$$IM = MI = M.$$

In this review, we consider matrix Lie-Santilli algebras and groups, with matrix product of a more general form than the conventional one. This generalized matrix product is based on the usual matrix product and also is associative and distributive. The corresponding unit matrix takes a more general form as compared to the standard one. Such algebras were first studied and developed by Santilli [3, 4]; see also recent reviews on fundamentals and applications of this theory in ref. [5].

In Sec. 2, we study admissible realizations of the axiom of matrix product [3]. Following Santilli, we conjecture that the most general form of the associative and distributive product of two $n \times n$ matrices M and N is of the form

$$M\hat{T}N,$$

where $\hat{T}$ is fixed $n \times n$ matrix, with the underlying product between M , $\hat{T}$, and N being the standard matrix one.

This form of product in algebras was first introduced and studied by Santilli [3], with $\hat{T}$ being called isotopic element. The axiom of unit element implies that the unit matrix $\hat{I}$, corresponding to this product, is of the form

$$\hat{I} = \hat{T}^{-1},$$

where the inverse matrix $\hat{T}^{-1}$ is defined due to the relationship $\hat{I}\hat{T} = \hat{T}\hat{I} = I$. Indeed, one can verify that

$$\hat{I}\hat{T}M = M\hat{T}\hat{I} = M$$

holds identically for any matrix M .

In the particular case when $\hat{I} = \hat{T}$, the standard unit matrix and standard matrix product are recovered, $\hat{I} = I = \hat{T}$. In the general case, $\hat{I} \neq \hat{T}$.

We introduce definition of *dual* matrix algebra, which is defined due to the interchange $\hat{I} \leftrightarrow \hat{T}$. Also, we consider metrics and associated coordinate systems in the matrix space. We study the conditions of reducing the unit matrix $\hat{I}$ to the standard one, and introduce the notion of generating matrix, which is of particular relevance for applications in physics.

Throughout the paper, we use the notation

$$M \hat{\times} N$$

adopted from ref. [6], to denote the associative and distributive matrix product between M and N . Matrix algebra of $n \times n$ matrices over field $\mathbb{C}$ equipped with an associative and distributive matrix product is denoted as $M(n, \mathbb{C}, \hat{\times})$.

In Sec. 3, we study some classical matrix Lie groups and Lie algebras, which are generalized to corresponding Lie-Santilli groups and Lie-Santilli algebras by assuming the general form of unit, and derive restrictions on the form of unit matrix $\hat{I}$ from original definitions of the Lie algebras and the associated Lie groups, i.e., obey the main principle of the theory stating that Lie character is preserved. By collecting all the derived restrictions we arrive at the conclusion that the most *simple* admissible nonstandard form of unit matrix appears to be of the form of a positive-definite diagonal matrix,

$$\hat{I} = \text{diag}(q_1, q_2, \dots, q_n),$$

$q_i > 0, i = 1, \dots, n$. This is in full accordance with earlier result by Santilli. We consider action of the Lie-Santilli groups on classical linear spaces $\mathbb{R}^n$ and $\mathbb{C}^n$, and find that the associated matrix Lie-Santilli groups, in which matrix product is $M\hat{T}N$ and unit is $\hat{I}$, conserve metrics which are, in general, *not conformally equivalent* to the Euclidean one. Also, we briefly review the infinite dimensional case, and indicate the relevance of the unit $\hat{I}$ in mathematical and physical context.

In Sec. 4, we consider in detail $SO(3, \mathbb{R}, \hat{\times})$, $SO(2, \mathbb{R}, \hat{\times})$, $SO(1, 1, \mathbb{R}, \hat{\times})$, and $U(1, \mathbb{C}, \hat{\times})$ matrix Lie-Santilli groups and $M(2, \mathbb{C}, \hat{\times})$ matrix Lie-Santilli algebra, as examples which are of interest in physics. To construct nontrivial realization (i.e., with the unit $\hat{I} \neq I$) of $SO(3, \mathbb{R}, \hat{\times})$, we use properties of the dual algebra. We establish the relationship between $SO(2, \mathbb{R}, \hat{\times})$ and $U(1, \mathbb{C}, \hat{\times})$, and consider the action of $U(1, \mathbb{C}, \hat{\times})$ on complex plane $\mathbb{C}$. It is remarkable to note that $U(1, \mathbb{C}, \hat{\times})$ makes, in general, linear *non complex-analytic* transformation of the complex plane $\mathbb{C}$. This is in confirmation of emphasize made by Santilli [6] that theory with the nonstandard form of product (and unit) is related to the standard one by a *non-unitary* transformation. We construct the realization map which relates $GL(n, \mathbb{C}, \hat{\times})$ with $GL(2n, \mathbb{R}, \hat{\times})$. One of the open problems is construction of a nontrivial realization of the $SU(2, \mathbb{C}, \hat{\times})$ Lie-Santilli group.

Detailed studies on associative algebras and groups with unit element different from the standard one, and its applications to physics, were made by Santilli [3, 4, 5, 6] since 1978; see also Sourlas and Tsagas [7]. We refer the reader to these papers for review and results of recent development of the Lie-Santilli algebras.

The main motivating idea lying behind the present talk, as well as some of the results, are due to the recent study presented in ref. [6]. Particularly, there it has been emphasized that the physical theories which are based on various types of classical and operator *deformations* of the standard ones should be

reformulated in order to provide their physical self-consistency and predictivity by preserving Lie character of the theories.

In the present talk, we consider different *realizations* of Santilli approach [6] trying to give self-consistent and detailed review of the matrix algebras.

Particularly, (i) we show that the associative and distributive product $M \hat{\times} N$ has a unique representation

$$M \hat{\times} N = M \hat{T} N,$$

as stated earlier by Santilli; (ii) introduce and use the so-called dual algebra $M(n, \mathbb{C}, \hat{\times}^{-1})$; (iii) analyze coordinate systems in the matrix space and homotopy class of unit; (iv) explicitly derive restrictions on the form of unit implied by self-consistent consideration of some classical Lie algebras; (v) explicitly establish the relationship between $SO(2, \mathbb{R}, \hat{\times})$ and $U(1, \mathbb{R}, \hat{\times})$, and the associated realization map for the case $n > 2$; and (vi) present nontrivial examples of some Lie-Santilli groups and algebras, with the specific examples of Santilli isounit.

Three important statements are in order.

i) Santilli [8] re-examined the classification of numbers and discovered that the abstract axioms of the numerical field $\mathbb{R}$ do not require necessarily that the basic unit be the usual number 1, but can be an arbitrary quantity subject that it is the inverse of T , and the new ring is equipped with product $M \hat{\times} N = M \hat{T} N$, where

$$I^*(t, r, v, a, E, \dots) = 1/T > 0$$

now universally known as Santilli isounit (viewed here as a scalar field over some variables t, r, v, a, E , etc.), T being known as the isotopic element, and 1 is the usual unit number. This led to the discovery of new numbers, those with an arbitrary positive definite unit, known as Santilli numbers (isoreal, isocomplex and isoquaternions) first presented in ref. [8]. After reformulating his Lie-isotopic theory on numbers, Santilli discovered that the emerging mathematical and physical theories remain inconsistent because, for instance, they are unable to predict the same numerical values under the same conditions at different times, as it is the case for Hamiltonian theories.

ii) Despite the above efforts, isotopic theories remained inconsistent. Following additional efforts, Santilli [9] found that the origin of the consistencies rests in the differential calculus that had been assumed for centuries not to depend on the basic unit. This led to the discovery of the differentiators calculus with the basic definition of the differential of the variable r

$$d^*r = Td[rI^*(t, r, v, \dots)],$$

where d stands for the usual differential.

iii) Reformulations of the realizations of Lie-Santilli groups and algebras, and representations of Lie-Santilli groups according to the isonumbers, isodifferential calculus, and isogeometry is a necessary task to avoid inconsistencies found by Santilli.

In fact, as one can see, the case 1×1 of Lie-Santilli matrix algebras over $\mathbb{R}$ represent natural generalization of the usual field of real numbers $\mathbb{R}$. This

requires formulation of field in consistency with Lie-Santilli algebras, and thus separate consideration of the construction of Lie-Santilli algebras over the field equipped with Santilli isounit (isofield). This has been done by Santilli [8, 9].

When dealing with the *realizations* of Lie-Santilli theory, which is the main issue of the present talk, we could not remain at the abstract axiom level and should choose specific isounit T for isofield, to provide actual calculations. The same is in the realization of Lie-Santilli matrix algebras where we should choose specific matrix isounit $\hat{I} \neq I = \text{diag}(1, 1, \dots, 1)$, for example, in the form $\hat{I} = \text{diag}(q_1, q_2, \dots, q_n)$, $q_i > 0$, $i = 1, \dots, n$, to study explicit examples of Lie-Santilli group representations.

While in the case of matrix algebras we can leave some parameters like q_i unspecified before making calculations, in the case of numbers there is no much room for such a parametric freedom, and we should specify T to be certain real number (for given values of the variables t, r, v, a , etc.), to make realization. Both usual field of real numbers $\mathbb{R}$ and any specific field of isonumbers (i.e., the field equipped with certain $T \neq 1$) are different *realizations* of the abstract field of real numbers. In the present talk we restrict ourselves by $\mathbb{R}$ (and $\mathbb{C}$) while the cases of isofields with different specific Santilli isounits are of much interest to consider since the general formalism necessitates consideration of mappings between them to avoid inconsistencies. One of the interesting issues is to consider Lie-Santilli algebras and groups modulo these mappings of the isofields.

For example, in higher dimensions, the consistency leads us to consider inhomogeneous dilation $\mathbb{R}_q^n$ (i.e., isospace) of the Euclidean space $\mathbb{R}^n$ (see Section 3.2.3).

2 Matrix algebra $M(n, \mathbb{C}, \hat{\times})$

2.1 The $\hat{\times}$ -product of matrices

In this subsection we present basics of Lie-Santilli matrix algebras.

We denote the set of all $n \times n$ matrices over field of complex numbers $\mathbb{C}$ by $M(n, \mathbb{C})$.

We start consideration by recalling that the usual matrix algebra $M(n, \mathbb{C})$, which is viewed as $M(n, \mathbb{C})$ equipped by the conventional matrix product, is a Lie algebra with respect to the commutator

$$[M, N] = MN - NM, \quad M, N \in M(n, \mathbb{C}), \quad (2.1)$$

where the product is usual product of matrices in the underlying associative and distributive algebra with standard unit $I = \text{diag}(1, 1, \dots, 1)$.

In $M(n, \mathbb{C})$, we follow Santilli [3] and define the $\hat{\times}$ -commutator as follows:

$$[M, N]_{\hat{\times}} = M\hat{\times}N - N\hat{\times}M, \quad M, N \in M(n, \mathbb{C}), \quad (2.2)$$

with respect to the $\hat{\times}$ -product

$$M\hat{\times}N = MTN, \quad (2.3)$$

i.e.

$$(M \hat{\times} N)_{ij} = \sum_{k,l=1}^n m_{ik} \hat{T}^{kl} n_{lj}, \quad (2.4)$$

where

$$\hat{T} = \hat{I}^{-1}, \quad \hat{T}\hat{I} = \hat{I}\hat{T} = I, \quad \hat{I} \in \mathcal{M}(n, \mathbb{C}), \quad (2.5)$$

and m_{ik} , $\hat{T}^{kl}$, and n_{lj} are corresponding matrix elements. Here, $\hat{I}$ is a *fixed* invertible matrix, which is the (left and right) unit, and we assume that in general $\hat{I} \neq \hat{T}$. Indeed, it can be immediately checked that $\hat{I}$ verifies the axiom of left and right unit element,

$$M \hat{\times} \hat{I} = \hat{I} \hat{\times} M = M, \quad \forall M \in \mathcal{M}(n, \mathbb{C}). \quad (2.6)$$

We denote the set $\mathcal{M}(n, \mathbb{C})$ equipped by the unit $\hat{I}$ and the associated $\hat{\times}$ -product as $\mathcal{M}(n, \mathbb{C}, \hat{\times})$.

Evidently, $\mathcal{M}(n, \mathbb{C}, \hat{\times})$ is a linear space, i.e., the following relationships hold:

$$M + N = N + M, \quad (2.7)$$

$$M + (N + P) = (M + N) + P, \quad (2.8)$$

$$M + 0 = M, \quad (2.9)$$

$$M + (-M) = 0, \quad (2.10)$$

$$\alpha(\beta M) = (\alpha\beta)M, \quad (2.11)$$

$$\alpha(M + N) = \alpha M + \alpha N, \quad (\alpha + \beta)M = \alpha M + \beta M, \quad (2.12)$$

$$1 \cdot M = M, \quad (2.13)$$

where α and β are complex numbers, and the $\hat{\times}$ -product is associative and distributive, i.e., the following two relationships hold (iso-associativity and iso-distributivity [3, 5]):

$$(M \hat{\times} N) \hat{\times} P = M \hat{\times} (N \hat{\times} P), \quad (2.14)$$

$$M \hat{\times} (N + P) = M \hat{\times} N + M \hat{\times} P, \quad (M + N) \hat{\times} P = M \hat{\times} P + N \hat{\times} P. \quad (2.15)$$

Inverse of matrix in the algebra $\mathcal{M}(n, \mathbb{C}, \hat{\times})$ is defined as

$$M^{-\hat{1}} \hat{\times} M = M \hat{\times} M^{-\hat{1}} = \hat{I}, \quad (2.16)$$

and multiplication of matrix by complex number α is as usual,

$$\alpha M = (\alpha m_{ij}). \quad (2.17)$$

This means that $\mathcal{M}(1, \mathbb{C}, \hat{\times})$ is assumed to be isomorphic to $\mathcal{M}(1, \mathbb{C})$. In one-dimensional case, the $\hat{\times}$ -product is $\alpha \hat{\times} \beta = \alpha \hat{T} \beta = \hat{T} \alpha \beta$, where α , β , and $\hat{T}$ are complex numbers, so that we can ignore overall fixed non-zero factor $\hat{T}$ in all the products. Indeed, there is an isomorphism between $\mathbb{C}$ and $\hat{T}\mathbb{C}$ provided by dilation. Nontriviality comes in *higher-dimensional* cases; see Secs. 3.2.3 and 4.2 for details.

We denote n th power of matrix in $M(n, \mathbb{C}, \hat{\times})$ by

$$M^{\hat{n}} = \underbrace{M \hat{\times} M \hat{\times} \cdots \hat{\times} M}_n, \quad (2.18)$$

and define $M^{\hat{0}} = \hat{I}$.

We conclude that $M(n, \mathbb{C}, \hat{\times})$ is an associative and distributive algebra.

Note that $M(n, \mathbb{C}, \hat{\times})$ preserves a Lie algebra character with respect to the $\hat{\times}$ -commutator. One can check that $\hat{\times}$ -commutator (2.2) is skew-symmetric, and Jacobi identity,

$$[M, [N, P]_{\hat{\times}}]_{\hat{\times}} + [P, [M, N]_{\hat{\times}}]_{\hat{\times}} + [N, [P, M]_{\hat{\times}}]_{\hat{\times}} = 0, \quad (2.19)$$

with respect to $\hat{\times}$ -commutator is satisfied. Namely,

$$\begin{aligned} [M, [N, P]_{\hat{\times}}]_{\hat{\times}} &= M \hat{\times} N \hat{\times} P - M \hat{\times} P \hat{\times} N - N \hat{\times} P \hat{\times} M + P \hat{\times} N \hat{\times} M, \\ [P, [M, N]_{\hat{\times}}]_{\hat{\times}} &= P \hat{\times} M \hat{\times} N - P \hat{\times} N \hat{\times} M - M \hat{\times} N \hat{\times} P + N \hat{\times} M \hat{\times} P, \\ [N, [P, M]_{\hat{\times}}]_{\hat{\times}} &= N \hat{\times} P \hat{\times} M - N \hat{\times} M \hat{\times} P - P \hat{\times} M \hat{\times} N + M \hat{\times} P \hat{\times} N, \end{aligned} \quad (2.20)$$

and by summing up these three expressions we obtain identically zero.

In definitions of usual Lie algebras $sl(n)$, $o(n)$, and $u(n)$, and associated Lie groups one uses the following operations with matrices: Trace, Transpose (M^t), and Complex Conjugate ($\bar{M}$). All these operations, and also Det, concern matrix elements and their definitions in the algebra $M(n, \mathbb{C}, \hat{\times})$ we keep the same as in algebra $M(n, \mathbb{C})$.

We emphasize that both algebras $M(n, \mathbb{C})$ and $M(n, \mathbb{C}, \hat{\times})$ obey the same set of axioms by construction.

Thus, the product (2.3) is one of the *admissible realizations* of abstract definition of the product in matrix algebra accompanied by associated realization of the unit element $\hat{I}$. This realization of product is based on the usual matrix product and can be thought of as the simplest generalization of it. However, note that $M(n, \mathbb{C}, \hat{\times})$ is not generalization of abstract matrix algebra. Instead, $M(n, \mathbb{C})$ and $M(n, \mathbb{C}, \hat{\times})$ are two *different realizations* of the abstract matrix algebra, with $M(n, \mathbb{C})$ being a *simplest realization* while $M(n, \mathbb{C}, \hat{\times})$ is an example of *more general realization* of it which is intimately based on $M(n, \mathbb{C})$.

Perhaps, some other admissible realizations of the matrix product exist. In general, this means that the axioms of matrix algebra do not fix the form of matrix product and the form of unit matrix to be only the standard ones. The aim of the present paper is to investigate implications of the form (2.3) of matrix product assuming that the matrix $\hat{I}$ is different from I .

As we will demonstrate in the following Sections, some severe restrictions on the form of unit matrix $\hat{I}$ naturally arise. Counterparts of the classical Lie groups associated to $M(n, \mathbb{C}, \hat{\times})$ reveal interesting properties. For example, the associated orthogonal group conserves the metrics defined by $\hat{I}$, which is, in general, *not conformally equivalent* to Euclidean metrics $I = (\delta_k^i)$. Also,

the associated unitary groups make, in general, linear *non complex-analytic* transformations of the complex space.

Following Santilli, we conjecture that the product (2.3) is the *most general* realization of the abstract associative and distributive product in matrix algebra with a unit. In Appendix A, we sketch the proof.

2.2 The dual algebra $M(n, \mathbb{C}, \hat{\times}^{-1})$

In this subsection we introduce the notion of dual algebra.

We start by noting that the transformation

$$\rho_0 : X \mapsto \hat{I}X\hat{I}, \quad X = M, N, \quad (2.21)$$

converts the $\hat{\times}$ -commutator

$$M\hat{T}N - N\hat{T}M \quad (2.22)$$

to the commutator with respect to $\hat{I}$,

$$M\hat{I}N - N\hat{I}M. \quad (2.23)$$

Indeed, we have

$$\begin{aligned} [\rho_0(M), \rho_0(N)]_{\hat{\times}} &= \hat{I}M\hat{I}\hat{T}\hat{I}N\hat{I} - \hat{I}N\hat{I}\hat{T}\hat{I}M\hat{I} \\ &= \hat{I}(M\hat{I}N - N\hat{I}M)\hat{I} = \rho_0(M\hat{I}N - N\hat{I}M) \neq \rho_0[M, N]_{\hat{\times}}. \end{aligned} \quad (2.24)$$

This means that ρ_0 is not endomorphism of Lie algebra $M(n, \mathbb{C}, \hat{\times})$. Vice versa, the transformation $\rho'_0 : X \mapsto \hat{T}X\hat{T}$ converts the commutator (2.23) to the $\hat{\times}$ -commutator.

In general, one can construct some algebra by making the following interchange:

$$\hat{I} \leftrightarrow \hat{T}. \quad (2.25)$$

We call the resulting algebra as a *dual* to $M(n, \mathbb{C}, \hat{\times})$, and denote it by $M(n, \mathbb{C}, \hat{\times}^{-1})$. In the dual algebra, $\hat{T}$ is a unit matrix while $\hat{I}$ is used in the definition of $\hat{\times}$ -product. Accordingly, we call commutator (2.23) as a *dual* commutator, which defines Lie algebra dual to the one specified by $\hat{\times}$ -commutator (2.2).

In the case $\hat{I} = I$, the two algebras $M(n, \mathbb{C}, \hat{\times})$ and $M(n, \mathbb{C}, \hat{\times}^{-1})$ degenerate to one algebra, $M(n, \mathbb{C})$, and this is the way to obey the *self-duality* condition, $M(n, \mathbb{C}, \hat{\times}) \simeq M(n, \mathbb{C}, \hat{\times}^{-1})$, which simply means that $\hat{I} = \hat{T}$. In other words, the usual matrix algebra $M(n, \mathbb{C})$ with standard form of unit is picked up by the self-duality condition. Indeed, $\hat{I} = \hat{T} \equiv \hat{I}^{-1}$ has only one nontrivial solution, $\hat{I} = \text{diag}(1, 1, \dots, 1)$, for the unit element.

Also, we note that the transformation

$$\rho_1 : X \mapsto \hat{I}X\hat{T} \equiv \hat{I}X\hat{I}^{-1} \quad (2.26)$$

is inner automorphism of $M(n, \mathbb{C}, \hat{\times})$ in terms of standard product. Indeed, it is homomorphism since

$$\rho_1(M) \hat{\times} \rho_1(N) = (\hat{I}M\hat{T})\hat{T}(\hat{I}N\hat{T}) = \hat{I}(M \hat{\times} N)\hat{T} = \rho_1(M \hat{\times} N), \quad (2.27)$$

maps $\hat{I}$ to $\hat{I}$, and makes one-to-one correspondence, with the inverse transformation

$$\rho_1^{-1} : X \mapsto \hat{T}X\hat{I}. \quad (2.28)$$

The transformations ρ_1 and ρ_1^{-1} make one-to-one correspondence between element X and its *conjugate*. It should be emphasized that $\hat{I}$ in Eq. (2.26) is fixed so that $\hat{I}X\hat{I}^{-1}$ do not, of course, form a class of conjugated elements for X . Instead, $\hat{I}X\hat{I}^{-1}$ form a *similarity* class ($\hat{I}$ is fixed and X is arbitrary).

Also, we note that the transformation

$$\rho_2 : X \mapsto S \hat{\times} X \hat{\times} S^{-1} \quad X, S \in M(n, \mathbb{C}, \hat{\times}), \quad (2.29)$$

is endomorphism of $M(n, \mathbb{C}, \hat{\times})$,

$$\begin{aligned} \rho_2(M) \hat{\times} \rho_2(N) &= (S \hat{\times} M \hat{\times} S^{-1}) \hat{T} (S \hat{\times} N \hat{\times} S^{-1}) \\ &= S \hat{\times} (M \hat{T} S^{-1} \hat{T} S \hat{T} N) \hat{\times} S^{-1} = S \hat{\times} (M \hat{\times} N) \hat{\times} S^{-1} = \rho_2(M \hat{\times} N), \end{aligned} \quad (2.30)$$

where we have used the definition (2.16).

2.3 Metrics and coordinates of $M(n, \mathbb{C}, \hat{\times})$

In this subsection we consider metrics and coordinate systems in space of matrices of Lie-Santilli algebra.

The $\hat{\times}$ -product (2.3) is a smooth function (polynomial) of matrix elements of multipliers M and N . We introduce metrics in $M(n, \mathbb{C}, \hat{\times})$ as follows:

$$|M|^2 = \sum_{i,j} |m_k^i| |\hat{T}_l^k| |m_j^l|. \quad (2.31)$$

Here, we have denoted $M = (m_k^i)$ and $\hat{T} = (\hat{T}_l^k)$, and naturally require $\hat{T}$ to be matrix of a *positive-definite* form. In the standard case [10], $\hat{T} = I \equiv (\delta_l^k)$, the metrics (2.31) is Euclidean, and simply computed as a sum of all squared matrix elements $|m_k^i|^2$, giving us in the result a real number; for example, $|I|^2 = n$ in $M(n, \mathbb{C})$.

For the metrics (2.31) we have, evidently,

$$|M + N| \leq |M| + |N|, \quad (2.32)$$

and also it can be easily proved that

$$|M \hat{\times} N| \leq |M| \hat{\times} |N|. \quad (2.33)$$

Indeed, we have for the inner product $\langle, \rangle$ with positive definite form $\hat{T}$

$$\langle x, \hat{T}x \rangle \langle y, \hat{T}y \rangle - \langle x, \hat{T}y \rangle^2 = \frac{1}{2}(\langle x, \hat{T}y \rangle - \langle y, \hat{T}x \rangle), \quad (2.34)$$

from which (2.33) follows.

Let us introduce local coordinates in the space of matrices, in the neighbourhood of $\hat{I}$,

$$|M - \hat{I}| \leq 1. \quad (2.35)$$

Coordinate $x(M)$ of matrix M in $M(n, \mathbb{C}, \hat{\times})$ is defined as

$$x_j^i(M) = m_j^i - \hat{I}_j^i, \quad x_j^i(\hat{I}) = 0. \quad (2.36)$$

If we multiply all the matrices, in the neighbourhood of $\hat{I}$, by the matrix $\hat{T} = \hat{I}^{-1}$ from the right then we can introduce the coordinate $y(M)$ of matrix M :

$$y_j^i(M) = m_j^k \hat{T}_k^i - \delta_j^i, \quad y_j^i(\hat{I}) = 0, \quad (2.37)$$

which is the coordinate of M in the neighbourhood of $I = (\delta_j^i)$. This coordinate system can be used for matrices M such that

$$|M - \hat{I}| \leq |\hat{I}|. \quad (2.38)$$

Thus, we have coordinate system $x(M)$ in the neighbourhood of $\hat{I}$, which is related to coordinate system $y(M)$ in the neighbourhood of I . The two coordinate systems coincide if $\hat{I} = I$.

In addition to the coordinate system $y(M)$, we can introduce the alternative one, with the multiplication from the left,

$$z_j^i(M) = \hat{T}_k^i m_j^k - \delta_j^i, \quad z_j^i(\hat{I}) = 0. \quad (2.39)$$

This coordinate system is *not* equivalent to $y(M)$ since in general $M\hat{T} \neq \hat{T}M$.

The following remarks are in order.

(a) The coordinate system $x(M)$ can be introduced in the neighbourhood of any matrix $\hat{I}$.

(b) The coordinate system $y(M)$ can be introduced for any *invertible* $\hat{I}$. The procedure of moving the neighbourhood, $x(M) \mapsto y(M)$, described above is formal and means that one can introduce coordinate system in the neighbourhood of any invertible element of the matrix space which is "equivalent" in some sense to the standard coordinate system in the neighbourhood of $I = (\delta_j^i)$.

(c) The choice of the center of coordinate system is a matter of convenience. Natural preference is made to the usual unit matrix $I = (\delta_j^i)$ as a center for the coordinate system. However, when one uses $\hat{I}$ as unit in the matrix algebra, as it is the case for $M(n, \mathbb{C}, \hat{\times})$, it becomes natural to choose $\hat{I}$ as a center of coordinate system to have a consistent picture. However, even in this case one can move the neighbourhood of $\hat{I}$ to the neighbourhood of $I = (\delta_j^i)$ by using coordinate system $y(M)$ or $z(M)$ because $\hat{I}$ is an invertible matrix, and $\hat{T}$ in

Eqs. (2.37) and (2.39) is always well defined. Note however that the center is still $\hat{I}$ due to the second equation in (2.37).

(d) It is remarkable to note that the sizes of the vicinities are different; see Eqs. (2.35) and (2.38).

(e) Also, we emphasize here that while the center of coordinate system in, e.g., Euclidean space $\mathbb{R}^n$ is indeed of no importance in accordance to its homogeneity, the choice of the center in matrix spaces, which are not in general homogeneous and commutative, is of some importance. This is reflected partially by the fact that we have two-fold way to rich neighbourhood of standard I , namely, coordinate systems $y(M)$ and $z(M)$.

2.4 Homotopy class of unit

In this subsection we study relationship between Santilli isounit $\hat{I}$ and usual unit I .

2.4.1 Path from $\hat{I}$ to I

Let us specify the form of unit $\hat{I}$ by picking up the diagonal form

$$\hat{I} = \text{diag}(q_1, q_2, \dots, q_n), \quad (2.40)$$

where parameters q_i satisfy the following conditions:

$$q_i \in \mathbb{R}, \quad q_i \neq 0 \quad (i = 1, 2, \dots, n), \quad \sum q_i \neq 0, \quad (2.41)$$

that is

$$\hat{I} \in M(n, \mathbb{R}), \quad \text{Det } \hat{I} \neq 0, \quad \text{Trace } \hat{I} \neq 0, \quad (2.42)$$

This specific form of $\hat{I}$ obeys the conditions (3.76) requiring that $\hat{I}$ should be real, symmetric, and non-traceless matrix, and the *diagonal* form of $\hat{I}$ appears to be important in universal definitions of algebras of pseudo-unitary and pseudo-orthogonal groups; see Sec. 3.3. Complete list of the requirements on the form of $\hat{I}$ will be presented in Sec. 3.4 below.

The matrix $\hat{T}$ is then given by

$$\hat{T} = \text{diag}(1/q_1, 1/q_2, \dots, 1/q_n). \quad (2.43)$$

The norm of the unit $\hat{I}$ defined by (2.31) is

$$|\hat{I}| = \sqrt{\sum q_i} = \sqrt{\text{Trace } \hat{I}}, \quad (2.44)$$

and not equal to zero due to (2.42). Then, to have real positive norm we must put

$$\text{Trace } \hat{I} > 0. \quad (2.45)$$

Below, we restrict consideration on the unit $\hat{I}$ of the form (2.40), which defines *n-parametric family of algebras* $M(n, \mathbb{C}, \hat{\times})$, with the parameters q_i .

We note that in the case

$$\hat{I} \rightarrow I, \quad (2.46)$$

if such a limit exists, we recover the original algebra $M(n, \mathbb{C})$. The limit does not always exist since $\hat{I}$ is a deformation of I by n real parameters $q_1, q_2, \dots, q_n$, which can have both negative and positive values, whereas $\hat{I}$ should always be invertible by definition. So, $\hat{I}$ and I should be *homotopically equivalent*, i.e. it must exist a smooth path in space $M(n, \mathbb{C})$ connecting $\hat{I}$ and I . This is possible if and only if q_i 's are positive numbers,

$$q_i > 0, \quad i = 1, 2, \dots, n. \quad (2.47)$$

Indeed, for negative value of some q_i , the path should go through the point $q_i = 0$, in which $\hat{I}$ is not an invertible matrix ($\text{Det } \hat{I} = 0$ at $q_i = 0$), and $\hat{I}$ blows up as $q_i \rightarrow 0$; see Eq. (2.43). Also, condition (2.47) follows from the requirement that $\hat{I}$ must be positive definite matrix; see Eq. (2.31).

Therefore, in general we must restrict consideration to the homotopy class of matrices to which standard unit matrix I belongs. Thus, the inverse map $I \mapsto \hat{I}$ is a smooth deformation, along with n -parametric path in space $M(n, \mathbb{C})$. The parameterization is given simply by the diagonal matrix $W = \text{diag}(w_1, w_2, \dots, w_n)$, with the parameters w_i running from 1 to q_i .

In view of the conditions (2.41) and (2.47), we can represent the matrix $\hat{I}$ as follows:

$$\hat{I} = \text{diag}(1 + r_1, 1 + r_2, \dots, 1 + r_n) \equiv I + R, \quad r_i > -1, \quad (2.48)$$

where $R = \text{diag}(r_1, r_2, \dots, r_n)$.

Whereas I commute with any element of $M(n, \mathbb{C})$, i.e. $[I, M] = 0$, one observes that $\hat{I}$ does not in general commute with any element of algebra $M(n, \mathbb{C})$,

$$[\hat{I}, M] \neq 0, \quad (2.49)$$

that is, $\hat{I}$ is not in the center of $M(n, \mathbb{C})$. However, it is in the center of the algebra $M(n, \mathbb{C}, \hat{\times})$,

$$[\hat{I}, M]_{\hat{\times}} = 0, \quad \forall M \in M(n, \mathbb{C}, \hat{\times}). \quad (2.50)$$

2.4.2 $\hat{I}$ in $M(2, \mathbb{C})$

To see more details on the connection between I and $\hat{I}$ and to provide an example, let us consider the usual algebra $M(2, \mathbb{C})$ of 2×2 complex matrices. The unit is

$$I = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix} \quad (2.51)$$

and we take the known example of matrix $\hat{I} \in M(2, \mathbb{C})$ of the form

$$\hat{I} = \begin{pmatrix} q_1 & 0 \\ 0 & q_2 \end{pmatrix}, \quad (2.52)$$

where

$$q_1 \neq 1, \quad q_2 \neq 1, \quad q_1 \neq q_2. \quad (2.53)$$

We emphasize here that I is unit matrix, and matrix product is an ordinary one; $\hat{I}$ is not unit matrix in $M(2, \mathbb{C})$.

We observe that both the matrices are *Hermitean*,

$$I^\dagger \equiv \bar{I}^t = I, \quad \hat{I}^\dagger = \hat{I}, \quad (2.54)$$

and therefore *normal*, i.e.,

$$I^\dagger I = II^\dagger, \quad \hat{I}^\dagger \hat{I} = \hat{I} \hat{I}^\dagger, \quad (2.55)$$

and therefore they are *simple*, i.e. multiplicity of each eigenvalue of the matrices is equal to its geometrical multiplicity.

While I is *unitary* matrix in $M(2, \mathbb{C})$,

$$I^\dagger I = I, \quad (2.56)$$

the matrix $\hat{I}$ is *not unitary* in $M(2, \mathbb{C})$,

$$\hat{I}^\dagger \hat{I} = \hat{I} \hat{I} = \text{diag}(q_1^2, q_2^2) \neq I. \quad (2.57)$$

It is remarkable to note, however, that in the algebra $M(2, \mathbb{C}, \hat{\times})$ we have

$$\hat{I}^\dagger \hat{\times} \hat{I} = \hat{I} \hat{\times} \hat{I} = \hat{I}, \quad (2.58)$$

so that $\hat{I}$ is unitary matrix in $M(2, \mathbb{C}, \hat{\times})$ while I is not unitary in $M(2, \mathbb{C}, \hat{\times})$.

Matrices I and $\hat{I}$ have different spectra and therefore they are not *unitary similar* to each other, i.e., there is no unitary matrix $U \in M(2, \mathbb{C})$ such that the relation

$$\hat{I} \neq U^{-1} I U \quad (2.59)$$

holds. Moreover, they are not even simply *similar* to each other, i.e.,

$$\hat{I} \neq S^{-1} I S, \quad (2.60)$$

for any matrix $S \in M(2, \mathbb{C})$. Indeed, $S^{-1} I S = S^{-1} S = I$ by definition, and thus this can not be equal to $\hat{I}$. So, vice versa,

$$I \neq V^{-1} \hat{I} V, \quad (2.61)$$

for any $V \in M(2, \mathbb{C})$ (see Appendix B for the proof).

The same properties are valid for higher dimensional cases, $n > 2$.

2.4.3 The generating matrix

Let us consider the ordinary eigenvalue problem,

$$(qI - Q)x = 0, \quad (2.62)$$

where $Q \in M(n, \mathbb{C})$. In the case Q is a positive-definite Hermitean matrix, we have set of positive real eigenvalues $q_1, q_2, \dots, q_n$, so that we can rewrite the above equation as

$$(\hat{I} - Q)x = 0, \quad (2.63)$$

where $\hat{I} = \text{diag}(q_1, q_2, \dots, q_n)$, $q_i > 0$. This means that Q and $\hat{I}$ are unitary similar to each other, $Q = U\hat{I}U^\dagger$, and $\hat{I}$ is positive-definite Hermitean matrix. We note that it is exactly the matrix we have as a general form of unit in $M(n, \mathbb{C}, \hat{\times})$; see also Sec. 3.4. Particularly, the scalar matrix $\hat{I} = \lambda I$ corresponds to fully degenerate spectrum of Q , in which case Q is necessarily of the form $Q = \lambda I$.

We call positive-definite Hermitean matrix Q satisfying Eq. (2.63) as a *generating matrix* for unit $\hat{I}$, $\hat{I} \in M(n, \mathbb{C}, \hat{\times})$. Clearly, all generating matrices for a given fixed $\hat{I}$ are unitary similar to each other, and they are not necessarily of a diagonal form.

If we relax the condition of diagonality of unit $\hat{I}$ (see Sec. 3.4) we can take generating matrix Q as a unit in algebra $M(n, \mathbb{C}, \hat{\times})$ provided that in some basis, unitary related to the original one, Q has a diagonal form.

3 Lie-Santilli groups and Lie-Santilli algebras

In this Section we study some classical Lie-Santilli groups and algebras which are based on the matrix algebra $M(n, \mathbb{C}, \hat{\times})$.

3.1 Lie-Santilli groups

3.1.1 Group $GL(n, \mathbb{C}, \hat{\times})$

We denote subgroup of $M(n, \mathbb{C}, \hat{\times})$ consisting of matrices obeying the condition of nonzero determinant,

$$\text{Det } M \neq 0, \quad M \in M(n, \mathbb{C}, \hat{\times}), \quad (3.1)$$

as $GL(n, \mathbb{C}, \hat{\times})$.

3.1.2 Unitary group $U(n, \mathbb{C}, \hat{\times})$

The group $U(n, \mathbb{C}, \hat{\times})$ is a subgroup of $GL(n, \mathbb{C}, \hat{\times})$ defined by the following unitarity condition:

$$U \hat{\times} \hat{I} \hat{\times} U^\dagger \equiv U \hat{\times} U^\dagger \equiv U \hat{T} U^\dagger = \hat{I}, \quad (3.2)$$

where we have denoted for Hermitean conjugation

$$U^\dagger \equiv \bar{U}^t. \quad (3.3)$$

From the unitarity condition (3.2), it follows that

$$\begin{aligned} \text{Det } (U \hat{\times} U^\dagger) &\equiv \text{Det } (U \hat{T} U^\dagger) = (\text{Det } \hat{T})(\text{Det } U)(\text{Det } U) \\ &= (\text{Det } \hat{T})|\text{Det } U|^2 = \text{Det } \hat{I}, \end{aligned} \quad (3.4)$$

that is,

$$|\text{Det } U|^2 = (\text{Det } \hat{I})^2, \quad U \in \text{U}(n, \mathbb{C}, \hat{\times}). \quad (3.5)$$

Apart from the usual case, determinant of unitary matrices in $\text{M}(n, \mathbb{C}, \hat{\times})$ is not, in general, equal to ± 1 .

We define the subgroup $\text{SU}(n, \mathbb{C}, \hat{\times})$ by the condition

$$|\text{Det } U| = \text{Det } \hat{I}, \quad U \in \text{SU}(n, \mathbb{C}, \hat{\times}), \quad (3.6)$$

i.e., determinant of special unitary matrices in $\text{M}(n, \mathbb{C}, \hat{\times})$ is equal to $\text{Det } \hat{I}$.

3.1.3 Orthogonal group $\text{O}(n, \mathbb{R}, \hat{\times})$

Similarly, for orthogonal group $\text{O}(n, \mathbb{R}, \hat{\times})$ we have

$$O \hat{\times} \hat{I} \hat{\times} O^t \equiv O \hat{\times} O^t \equiv O \hat{T} O^t = \hat{I}, \quad O \in \text{O}(n, \mathbb{R}, \hat{\times}), \quad (3.7)$$

and for $\text{SO}(n, \mathbb{R}, \hat{\times})$ we have additionally,

$$\text{Det } O = \text{Det } \hat{I}. \quad (3.8)$$

3.1.4 Group $\text{SL}(n, \mathbb{C}, \hat{\times})$

The subgroup $\text{SL}(n, \mathbb{C}, \hat{\times})$ of $\text{GL}(n, \mathbb{C}, \hat{\times})$ is defined due to condition

$$\text{Det } M = \text{Det } \hat{I}, \quad M \in \text{SL}(n, \mathbb{C}, \hat{\times}). \quad (3.9)$$

Note that for the usual case, $\hat{I} = I$, we have $\text{Det } \hat{I} = 1$ while for $\hat{I}$ of the form (2.40) we have

$$\text{Det } \hat{I} = q_1 q_2 \cdots q_n. \quad (3.10)$$

3.2 Action of the groups on classical linear spaces

3.2.1 Eigenvalue problem

We define natural (left) action of the group $\text{GL}(n, \mathbb{C}, \hat{\times})$ on complex space $\mathbb{C}^n$ as

$$z \mapsto M \hat{\times} z = M \hat{T} z, \quad M \in \text{GL}(n, \mathbb{C}, \hat{\times}), \quad z \in \mathbb{C}^n. \quad (3.11)$$

This definition is consistent with the action of unit $\hat{I}$,

$$\hat{I} \hat{\times} z = z, \quad (3.12)$$

which can be viewed as an identity transformation of $\mathbb{C}^n$. Also, for two consequential actions, we have

$$z' = M \hat{\times} z, \quad z'' = N \hat{\times} z' = N \hat{\times} (M \hat{\times} z) = (N \hat{\times} M) \hat{\times} z = Q \hat{\times} z, \quad (3.13)$$

$$M, N, Q \in \text{GL}(n, \mathbb{C}, \hat{\times}),$$

that means that this action is consistent with the algebra $\text{GL}(n, \mathbb{C}, \hat{\times})$.

The eigenvalue problem is then defined by the following equation:

$$M \hat{\times} z = \lambda z, \quad (3.14)$$

where $\lambda \in \mathbb{C}$, or, equivalently,

$$(\lambda \hat{I} - M) \hat{\times} z = 0. \quad (3.15)$$

This equation can be identically rewritten as

$$(\lambda I - M \hat{T}) z = 0. \quad (3.16)$$

Thus, in the algebra $M(n, \mathbb{C}, \hat{\times})$ the characteristic polynomial of the matrix M is

$$c(\lambda) = \text{Det} (\lambda I - M \hat{T}). \quad (3.17)$$

3.2.2 Action of unitary group

Let us identify (Hermitean) scalar product in $\mathbb{C}^n$,

$$\langle z_1, z_2 \rangle_{\mathbb{C}} = \sum_{i,j=1}^n z_1^i g_{ij} \bar{z}_2^j, \quad (3.18)$$

where $z_{1,2} \in \mathbb{C}^n$, such that it is conserved under the action of unitary matrix $U \in \text{U}(n, \mathbb{C}, \hat{\times})$ on each multiplier,

$$\langle U \hat{\times} z_1, U \hat{\times} z_2 \rangle = \langle z_1, z_2 \rangle. \quad (3.19)$$

From the above two equations we have immediately

$$\sum U_k^i \hat{T}_m^k z^m g_{ij} \bar{z}^n \hat{T}_n^l \bar{U}_l^j = \sum z^m g_{mn} \bar{z}^n, \quad (3.20)$$

i.e.

$$U_k^i \hat{T}_m^k g_{ij} \hat{T}_n^l \bar{U}_l^j = g_{mn}, \quad (3.21)$$

or

$$U \hat{\times} g \hat{\times} U^\dagger = g. \quad (3.22)$$

Therefore, we should put $g = \hat{I}$ to achieve consistency with the definition of unitarity (3.2). Obviously, this does not mean that $\hat{I}$ is conserved by the unitary matrix in the sense $U \hat{I} U^\dagger = \hat{I}$. Instead, the matrix

$$\tilde{U} = U \hat{T}, \quad (3.23)$$

with $\tilde{U}^\dagger = \hat{T}U^\dagger$, plays such a role, namely, we have from Eq. (3.22)

$$\tilde{U}\hat{I}\tilde{U}^\dagger = \hat{I}, \quad (3.24)$$

which is consistent with the definition (3.2).

Note that the scalar product (3.18), with $g = \hat{I}$, is Hermitean since $\hat{I}$ is real symmetric positive-definite matrix (Hermitean, $\hat{I}^\dagger = I$).

3.2.3 Action of orthogonal group

Similarly to the above case, the group $O(n, \mathbb{C}, \hat{\times})$ conserves the following scalar product:

$$\langle x_1, x_2 \rangle_{\mathbb{R}} = \sum_{i,j=1}^n \hat{I}_{ij} x_1^i x_2^j, \quad (3.25)$$

where $x_{1,2} \in \mathbb{R}^n$, or explicitly,

$$\langle x_1, x_2 \rangle_{\mathbb{R}} = q_1 x_1^1 x_2^1 + q_2 x_1^2 x_2^2 + \cdots + q_n x_1^n x_2^n. \quad (3.26)$$

This scalar product defines Euclidean space $\mathbb{R}^n$ equipped by the metrics $\hat{I}$,

$$\langle x, x \rangle = \sum_{i=1}^n q_i (x^i)^2, \quad (3.27)$$

i.e. the metric tensor is

$$g_{ij} = q_i \delta_{ij}. \quad (3.28)$$

We denote Euclidean space $\mathbb{R}^n$ equipped by the metrics (3.28) by $\mathbb{R}_q^n$. The matrix conserving the above scalar product is of the form

$$\tilde{O} = O\hat{T}, \quad (3.29)$$

i.e.

$$\tilde{O}\hat{I}\tilde{O}^t = \hat{I}. \quad (3.30)$$

Note that

$$\langle z, z \rangle_{\mathbb{C}} = \langle x, x \rangle_{\mathbb{R}}, \quad (3.31)$$

as in the usual case.

In terms of the algebra $M(n, \mathbb{C}, \hat{\times})$, the above unitarity and orthogonality definitions mean that the matrix $\hat{I}$ is an invariant. In contrast, in terms of the usual algebra $M(n, \mathbb{C})$ and geometry, these mean that the metric tensor $\hat{T}$ is transformed to metric tensor $\hat{I}$; see Eqs. (3.2) and (3.7). In the limiting case $\hat{I} = I = \hat{T}$, we have conservation of the metric tensor δ_{ij} .

This situation can be readily understood in terms of the duality property (2.25) of algebra $M(n, \mathbb{C}, \hat{\times})$. Namely, the definitions relate *dual spaces*, the one equipped by the metric (3.28) and the other equipped by the metric $\hat{T}$, i.e.,

$$g_{ij}^{dual} = \frac{1}{q_i} \delta_{ij}, \quad (3.32)$$

which are simply inverse to each other, and coincide when $\hat{I} = I = \hat{T}$.

Note that the above space $\mathbb{R}_q^n$ with metrics (3.27) can not be obtained from Euclidean space $\mathbb{R}^n$ by dilation $x \mapsto \lambda x$, except for one-dimensional case. Instead, we have the transformation

$$x^i \mapsto \hat{x}^i = x^i / \sqrt{q_i}, \quad (3.33)$$

which we call *inhomogeneous dilation*, $\hat{x} \in \mathbb{R}_q^n$. So the map $I \mapsto \hat{I}$ does not correspond in general to any linear conformal transformation of $\mathbb{R}^n$. Only when $q_1 = q_2 = \dots = q_n$ this is the case. The transformation (3.33) can be thought of as that it gives the coordinates x^i different *weights*.

Accordingly, the use of the general form (2.40) of unit $\hat{I}$ assumes, in general, *different weights* of the coordinates in contrast to *equal weights* provided by the standard unit I .

The equations $g_{ij}x^ix^j = (\text{Det } \hat{I})^2$ and $g_{ij}^{dual}x^ix^j = (\text{Det } \hat{T})^2$ define *fundamental ellipsoids*,

$$q_1(x^1)^2 + q_2(x^2)^2 + \dots + q_n(x^n)^2 = (q_1q_2 \dots q_n)^2, \quad (3.34)$$

$$\frac{1}{q_1}(x^1)^2 + \frac{1}{q_2}(x^2)^2 + \dots + \frac{1}{q_n}(x^n)^2 = \frac{1}{(q_1q_2 \dots q_n)^2}, \quad (3.35)$$

corresponding to $\hat{I}$ and $\hat{T}$, respectively. They are regular $(n-1)$ -hypersurfaces in $\mathbb{R}^n$. In usual terms, definitions of the unitarity and orthogonality are such that corresponding matrices transform the second ellipsoid to the first one. The sphere $\sum (x^i)^2 = 1$ lies between the ellipsoids, and is a limiting case of both the ellipsoids.

The following remarks are in order.

(a) In usual geometrical terms, these ellipsoids are not conserving under the orthogonality group $O(n, \mathbb{C}, \hat{\times})$. So, none of which is a homogeneous space of this group, and group $O(n, \mathbb{C}, \hat{\times})$ does not act on it transitively in the usual sense. Indeed, varying matrix O in Eq. (3.7), we observe that they act on $\hat{T}\hat{I}\hat{T} = \hat{T}$, which is a fixed matrix, and the result is another fixed matrix $\hat{I}$. However, in terms of the group $O(n, \mathbb{C}, \hat{\times})$, the ellipsoid (3.34) defined by $\hat{I}$ is *conserved* due to Eq. (3.7). So this ellipsoid is a homogeneous space of group $O(n, \mathbb{C}, \hat{\times})$ under the action of this group on it, and every two points of the ellipsoid can be connected by some $O \in O(n, \mathbb{C}, \hat{\times})$ (transitivity).

(b) Also, we see from the above considerations that the matrix of the form

$$\tilde{M} = M\hat{T} \quad (3.36)$$

appears to be of frequent use in the formalism. Note that according to Eq. (2.49), we have in general $M\hat{T} \neq \hat{T}M$. We shall see in Sec. 3.3 that matrices of the form $M\hat{T}$ is also of use in the Lie algebras. From Eq. (2.37), we see that such matrices correspond to those described in the neighbourhood of standard unit $I = (\delta_{ij})$.

(c) The use of different possible definitions of unitarity (3.2) and orthogonality (3.7)

$$U^\dagger \hat{T} U = \hat{I}, \quad O^t \hat{T} O = \hat{I}, \quad (3.37)$$

respectively, yields the same set up as above, with the matrix $M\hat{T}$ replaced by $\hat{T}M$. Note that this corresponds to choosing of the coordinate system (2.39) instead of (2.37).

3.2.4 Action of pseudo-unitary and pseudo-orthogonal groups

Definitions given in Sec. 3.1 can be extended to the case of pseudo-Euclidean spaces, with accordingly defined pseudo-unitary group $U(m, k, \mathbb{C}, \hat{\times})$ and pseudo-orthogonal group $O(m, k, \mathbb{R}, \hat{\times})$.

Let us define the metrics

$$\hat{G}_r = G\hat{I}, \quad (3.38)$$

where

$$G = \text{diag}(\underbrace{1, 1, \dots, 1}_m, \underbrace{-1, -1, \dots, -1}_k) \quad (3.39)$$

is metrics of pseudo-Euclidean space $\mathbb{R}^{m,k}$. Then, definitions of pseudo-unitary group $U_r(m, k, \mathbb{C}, \hat{\times})$ and pseudo-orthogonal group $O_r(m, k, \mathbb{R}, \hat{\times})$ are, respectively,

$$U \hat{\times} \hat{G}_r \hat{\times} U^\dagger = \hat{G}_r, \quad O \hat{\times} \hat{G}_r \hat{\times} O^t = \hat{G}_r. \quad (3.40)$$

These groups conserve the metrics $\hat{G}_r$. The other possible definition of metrics,

$$\hat{G}_l = \hat{I}G, \quad (3.41)$$

leads to definitions of different groups $U_l(m, k, \mathbb{C}, \hat{\times})$ and $O_l(m, k, \mathbb{R}, \hat{\times})$,

$$U \hat{\times} \hat{G}_l \hat{\times} U^\dagger = \hat{G}_l, \quad O \hat{\times} \hat{G}_l \hat{\times} O^t = \hat{G}_l, \quad (3.42)$$

since in general $G\hat{I} \neq \hat{I}G$. Also, note that

$$[\hat{G}_r, \hat{G}_l] = G\hat{I}\hat{I}G - \hat{I}G\hat{I}G \neq 0 \quad (3.43)$$

and

$$[\hat{G}_r, \hat{G}_l]_{\hat{\times}} = G\hat{I}G - \hat{I}G\hat{I}G \neq 0. \quad (3.44)$$

The groups $U_l(m, k, \mathbb{C}, \hat{\times})$ and $O_l(m, k, \mathbb{R}, \hat{\times})$ conserve the metrics $\hat{G}_l$. Evidently, these definitions of the groups are directly equivalent to that with respect to (3.38) if and only if pseudo-Euclidean metrics G and unit $\hat{I}$ commute, $G\hat{I} = \hat{I}G$, so that $\hat{G}_r$ and $\hat{G}_l$ coincide,

$$\hat{G}_r = \hat{G}_l = \hat{G}, \quad (3.45)$$

and we can put

$$U \hat{\times} \hat{G} \hat{\times} U^\dagger = \hat{G}, \quad O \hat{\times} \hat{G} \hat{\times} O^t = \hat{G}, \quad (3.46)$$

for definitions of pseudo-unitary group $U(m, k, \mathbb{C}, \hat{\times})$ and pseudo-orthogonal group $O(m, k, \mathbb{R}, \hat{\times})$, respectively. This is the case only for *diagonal* form of the unit $\hat{I}$ because in general only diagonal matrices commute with pseudo-Euclidean metrics G .

Note that due to the inner automorphism (2.26) we have $\rho_1 : \hat{G}_r \mapsto \hat{G}_l$, namely, $\hat{G}_l = \hat{I}\hat{G}_r\hat{I}$, so that $\hat{G}_r$ and $\hat{G}_l$ are elements conjugated to each other in $\text{GL}(n, \mathbb{C}, \hat{\times})$, and the groups $U_r(m, k, \mathbb{C}, \hat{\times})$ ($O_r(m, k, \mathbb{R}, \hat{\times})$) and $U_l(m, k, \mathbb{C}, \hat{\times})$ ($O_l(m, k, \mathbb{R}, \hat{\times})$) are conjugated to each other, as subgroups of $\text{GL}(n, \mathbb{C}, \hat{\times})$.

3.3 Matrix exponent and Lie-Santilli algebras

In this Section we study realizations and relationship between matrix exponents and Lie-Santilli algebras.

Tangent spaces in the neighbourhood of the unit $\hat{I}$ for the above considered groups are corresponding Lie algebras $gl(n, \mathbb{C}, \hat{\times})$, $sl(n, \mathbb{C}, \hat{\times})$, $u(n, \mathbb{C}, \hat{\times})$, and $o(n, \mathbb{R}, \hat{\times})$, which are well defined Lie algebras, as in the usual case.

The map from the tangent spaces to the groups is achieved by matrix exponent. The matrix exponent is defined, as usually, due to its formal series expansion. In $M(n, \mathbb{C}, \hat{\times})$, we define

$$\hat{e}^M = \sum_{n=0}^{\infty} \frac{M^{\hat{n}}}{n!}, \quad (3.47)$$

where the $\hat{n}$ -power of matrix M is defined by Eq. (2.18), and we put

$$\hat{e}^0 = \hat{I}. \quad (3.48)$$

Explicitly,

$$\hat{e}^M = \hat{I} + M + \frac{1}{2!} M \hat{\times} M + \dots. \quad (3.49)$$

This series expansion converges due to Eqs. (2.32) and (2.33). Then, one can easily prove using Eqs. (3.47) and (3.48) that

$$\hat{e}^{M+N} = \hat{e}^M \hat{\times} \hat{e}^N, \quad \text{for } \hat{\times}\text{-commuting matrices } M \text{ and } N, \quad (3.50)$$

$$\text{If } M = \hat{e}^X, \text{ then exists } M^{-1} = \hat{e}^{-X}, \quad (3.51)$$

$$\hat{e}^{X^t} = (\hat{e}^X)^t. \quad (3.52)$$

The above definition of matrix exponent defines *local coordinates* in the tangent space of group in the neighbourhood of unit element $\hat{I}$ of the group which have the following explicit form:

$$x_j^i(M) = (\ln M)_j^i = (M - \hat{I})_j^i - \dots, \quad (3.53)$$

where M is a group element. This map is one-to-one correspondence in some neighbourhood of the point $x_j^i(M) = 0$.

The matrix exponent in $M(n, \mathbb{C}, \hat{\times})$ is simply related to the usual matrix exponent by

$$\hat{e}^M = \hat{I} e^{\hat{T}M}, \quad M \in \text{GL}(n, \mathbb{C}, \hat{\times}), \quad (3.54)$$

with $e^0 = I$. Indeed, by using the power expansion we have

$$\hat{e}^M = \hat{I} + M + \frac{1}{2!} M \hat{T} M + \dots = \hat{I} (I + \hat{T} M + \frac{1}{2!} \hat{T} M \hat{T} M + \dots) = \hat{I} e^{\hat{T}M}. \quad (3.55)$$

The following remarks are in order.

(i) In fact, we need only in the neighbourhood of unit $\hat{I}$ when dealing with Lie algebras. In general, matrix exponent is not one-to-one correspondence when it is extended to the whole group; well known example is the usual $SL(2, \mathbb{R})$.

(ii) Note that there is an alternative relation,

$$\hat{e}^M = e^{M\hat{T}} \hat{I}, \quad (3.56)$$

between the matrix exponents. This relation is equivalent to (3.55), in the algebra $M(n, \mathbb{C}, \hat{\times})$. Indeed, let us check that the r.h.s. of (3.54) $\hat{\times}$ -commute with the r.h.s. of (3.56), in the neighbourhood of unit,

$$\begin{aligned} [e^{M\hat{T}} \hat{I}, \hat{I}e^{\hat{T}M}]_{\hat{\times}} &= e^{M\hat{T}} \hat{I}e^{\hat{T}M} - \hat{I}e^{\hat{T}M} \hat{T}e^{M\hat{T}} \hat{I} \\ &\simeq (I + M\hat{T})\hat{I}(I + \hat{T}M) - \hat{I}(I + \hat{T}M)\hat{T}(I + M\hat{T})\hat{I} \\ &= (I + 2M + M\hat{T}M) - (I + 2M + M\hat{T}M) = 0, \end{aligned} \quad (3.57)$$

where we have dropped higher-order terms.

Using the matrix exponent (3.47) we can prove that if $M \in SL(n, \mathbb{C}, \hat{\times})$, i.e. $\text{Det } M = \text{Det } \hat{I}$, then algebra $sl(n, \mathbb{C}, \hat{\times})$, as a tangent space of the group in the neighbourhood of unit $\hat{I}$, consists of matrices X such that

$$\text{Trace } X\hat{T} = 0, \quad (3.58)$$

and vice versa.

Indeed, let $\text{Trace } X\hat{T} = 0$. For $M(t) = \hat{e}^{tX}$ we have

$$M(t_1 + t_2) = M(t_1) \hat{\times} M(t_2), \quad (3.59)$$

where we used the fact that $M(t_1) = \hat{e}^{t_1 X}$ and $M(t_2) = \hat{e}^{t_2 X}$ $\hat{\times}$ -commute. Therefore,

$$\text{Det } M(t_1 + t_2) = \text{Det } M(t_1) \text{Det } \hat{T} \text{Det } M(t_2). \quad (3.60)$$

Solution for this equation is given by $F(t) \equiv \text{Det } M(t) = c_1 e^{c_2 t}$, where $c_{1,2}$ are constants. Evidently, $c_1 = (\text{Det } \hat{T})^{-1} = \text{Det } \hat{I}$. On the other hand,

$$\begin{aligned} F(t) &= \text{Det } \hat{e}^{tX} = \text{Det } (\hat{I} + tX + o(t)) = \text{Det } (I + t\hat{T}X + o(t))\hat{I} \\ &= (t \text{Trace } X\hat{T} + o(t))(\text{Det } \hat{I}). \end{aligned} \quad (3.61)$$

So, if $\text{Trace } X\hat{T} = 0$ then

$$c_2 = \frac{1}{c_1} \frac{dF}{dt} \Big|_{t=0} = \text{Trace } X\hat{T} = 0. \quad (3.62)$$

Thus we have, finally, $F(t) = \text{Det } \hat{I}$, i.e. $\text{Det } M = \text{Det } \hat{I}$. It can be easily shown that, vice versa, if $\text{Det } M = \text{Det } \hat{I}$ then $\text{Trace } X\hat{T} = 0$.

Comparing this to the usual relation, $\text{Trace } X = 0$, for $sl(n, \mathbb{C})$ we see some modification. Let us denote

$$\hat{\text{Trace}} M = \text{Trace } M\hat{T}. \quad (3.63)$$

One can verify that Trace M is not conserved under unitary transformation while $\hat{\text{Trace}} M$ is a conserved entity. Indeed, let us make the unitary transformation

$$M \mapsto M' = U \hat{\times} M \hat{\times} U^\dagger, \quad (3.64)$$

where unitary matrix U is given by (3.2). Then,

$$\begin{aligned} \hat{\text{Trace}} M' &= \sum M'_{ij} \hat{T}^{ji} = \sum (U_{ik} \hat{T}^{km} M_{mn} \hat{T}^{nl} U_{lj}^\dagger) \hat{T}^{ji} \\ &= \sum \hat{I}_{kl} \hat{T}^{km} M_{mn} \hat{T}^{nl} = \sum \delta_l^m M_{mn} \hat{T}^{nl} = \sum M_{mn} \hat{T}^{nm} = \hat{\text{Trace}} M, \end{aligned} \quad (3.65)$$

where we used the unitarity condition $\sum U_{lj}^\dagger \hat{T}^{ji} U_{ik} = \hat{I}_{lk}$, and the fact that $\sum \hat{I}_{kl} \hat{T}^{km} = \sum \hat{I}_{lk} \hat{T}^{km} = \delta_l^m$.

Below, we investigate explicitly relations between the groups $U(n, \mathbb{C}, \hat{\times})$, $O(n, \mathbb{R}, \hat{\times})$ and their tangent spaces.

Let us consider neighbourhood of the unit $\hat{I}$ of the group $U(n, \mathbb{C}, \hat{\times})$. Let $U(t) \in U(n, \mathbb{C}, \hat{\times})$ and $U(0) = \hat{I}$, where t is parameter. Then, we have

$$U(t) \hat{\times} U^\dagger(t) = \hat{I}, \quad \left. \frac{dU}{dt} \right|_{t=0} = X, \quad (3.66)$$

where X belongs to tangent space of $U(n, \mathbb{C}, \hat{\times})$ in the neighbourhood of $\hat{I}$. Differentiating first equation in (3.66), we have

$$\begin{aligned} \left. \frac{d}{dt} (U \hat{T} U^\dagger) \right|_{t=0} &= \left[\left. \frac{dU}{dt} \hat{T} U^\dagger + U \hat{T} \frac{dU^\dagger}{dt} \right] \right|_{t=0} \\ &= X \hat{T} \hat{I} + \hat{I} \hat{T} X^\dagger = X + X^\dagger = 0. \end{aligned} \quad (3.67)$$

So, $u(n, \mathbb{C}, \hat{\times})$ consists of the skew-Hermitean matrices,

$$X = -X^\dagger. \quad (3.68)$$

Similarly, it can be shown that if $O \in O(n, \mathbb{R}, \hat{\times})$ is orthogonal matrix then matrices from the tangent space in the neighbourhood of unit of this group are skew-symmetric, and vice versa. Indeed, for X such that

$$X = -X^t, \quad (3.69)$$

we have

$$O \hat{\times} O^t = \hat{e}^X \hat{\times} (\hat{e}^X)^t = \hat{e}^X \hat{\times} \hat{e}^{X^t} = \hat{e}^{X+X^t} = \hat{I}, \quad (3.70)$$

where we used Eqs. (3.50) and (3.52), and the fact that X and X^t are $\hat{\times}$ -commuting matrices.

For the tangent space elements X of pseudo-unitary group $U(m, k, \mathbb{C}, \hat{\times})$ and pseudo-orthogonal group $O(m, k, \mathbb{C}, \hat{\times})$, it is an easy exercise to obtain from Eqs. (3.46) the usual relations,

$$XG + GX^\dagger = 0, \quad XG + GX^t = 0, \quad (3.71)$$

respectively, where G is matrix of pseudo-Euclidean metrics (3.39). For example, for the pseudo-unitary group $U(m, k, \mathbb{C}, \hat{\times})$ we have

$$\begin{aligned} \frac{d}{dt}(U\hat{T}\hat{G}\hat{T}U^\dagger)|_{t=0} &= \left[\frac{dU}{dt}\hat{T}\hat{G}\hat{T}U^\dagger + U\hat{T}\hat{G}\hat{T}\frac{dU^\dagger}{dt} \right]_{t=0} \\ &= X(\hat{T}\hat{G}\hat{T})\hat{I} + \hat{I}(\hat{T}\hat{G}\hat{T})X^\dagger = XG + GX^\dagger = 0, \end{aligned} \quad (3.72)$$

where we have assumed that the matrices G and $\hat{I}$ commute and thus $\hat{I}$ is *diagonal* matrix; see remark below Eq. (3.46). In general, we have instead of (3.71),

$$X\hat{T}\hat{G}\hat{I} + GX^\dagger = 0, \quad X\hat{T}\hat{G}\hat{I} + GX^t = 0, \quad (3.73)$$

for the groups $U_r(m, k, \mathbb{C}, \hat{\times})$ and $O_r(m, k, \mathbb{R}, \hat{\times})$, respectively, where $\hat{G} = G\hat{I} \equiv \hat{G}_r$, and

$$XG + \hat{I}G\hat{T}X^\dagger = 0, \quad XG + \hat{I}G\hat{T}X^t = 0, \quad (3.74)$$

for the groups $U_l(m, k, \mathbb{C}, \hat{\times})$ and $O_l(m, k, \mathbb{R}, \hat{\times})$, respectively, where $\hat{G} = \hat{I}G \equiv \hat{G}_l$. The definitions (3.73) and (3.74) can be rewritten in a compact natural form,

$$X\hat{\times}\hat{G} + \hat{G}\hat{\times}X^\dagger = 0, \quad X\hat{\times}\hat{G} + \hat{G}\hat{\times}X^t = 0, \quad (3.75)$$

where $\hat{G} = \hat{G}_r$, or $\hat{G} = \hat{G}_l$.

To prove that the above tangent spaces indeed are Lie algebras one must show that the following properties hold:

- 1) If $\text{Trace } M = 0$ and $\text{Trace } N = 0$ then $\text{Trace } [M, N]_{\hat{\times}} = 0$.
- 2) If $M^t = -M$ and $N^t = -N$ then $[M, N]_{\hat{\times}}^t = -[M, N]_{\hat{\times}}$.
- 3) If $M^\dagger = -M$ and $N^\dagger = -N$ then $[M, N]_{\hat{\times}}^\dagger = -[M, N]_{\hat{\times}}$.

In fact, by this one shows that the spaces $sl(n, \mathbb{C}, \hat{\times})$, $o(n, \mathbb{R}, \hat{\times})$, and $u(n, \mathbb{C}, \hat{\times})$ are closed with respect to $\hat{\times}$ -commutator.

Note that the unit $\hat{I}$ is subject to Trace, Transpose and Complex conjugate operations in the above 1)-3). Let us put the following restrictions on $\hat{I}$:

$$\text{Trace } \hat{I} \neq 0, \quad \hat{I}^t = \hat{I}, \quad \bar{\hat{I}} = \hat{I}, \quad (3.76)$$

i.e. $\hat{I}$ is real symmetric matrix with non-zero trace. Particularly, the form (2.40) obeys the requirements (3.76).

Then, it is easy to check that the properties 1)-3) hold for any $\hat{I}$ obeying (3.76). Namely, we have

$$1) \quad \text{Trace } [M, N]_{\hat{\times}} = \text{Trace } M\hat{T}N - \text{Trace } N\hat{T}M \quad (3.77)$$

$$= \text{Trace } M\hat{T}N - \text{Trace } M\hat{T}N = 0.$$

$$\begin{aligned} 2) \quad [M, N]_{\hat{\times}}^t &= (M\hat{T}N)^t - (N\hat{T}M)^t = N^t\hat{T}^tM^t - M^t\hat{T}^tN^t \\ &= N^t\hat{T}M^t - M^t\hat{T}N^t = N\hat{T}M - M\hat{T}N \\ &= -[M, N]_{\hat{\times}}. \end{aligned} \quad (3.78)$$

Similarly, for 3).

In the same manner, one can prove that tangent spaces of $U(m, k, \mathbb{C}, \hat{\times})$ and $O(m, k, \mathbb{R}, \hat{\times})$ are Lie algebras with respect to $\hat{\times}$ -commutator.

3.4 Some restrictions on the form of Santilli isounit

We are now in a position to collect all the restrictions on the form of $\hat{I}$ stemming from consideration made above.

1) To have well defined inverse required to set up $\hat{\times}$ -product, $\hat{I}$ should be *non-degenerate*; see Sec. 2.1.

2) To define positive norm in space of matrices, $\hat{I}$ should be matrix of a *positive-definite symmetric* form; see Sec. 2.3.

3) To have positive norm, $\hat{I}$ should have *positive trace*; see Sec. 2.3.

4) To be in the homotopy class of I , $\hat{I}$ should have *positive* diagonal elements, for a diagonal form of $\hat{I}$; see Sec. 2.4.

5) For consistent definitions of algebras of orthogonal and unitary groups, $\hat{I}$ should be *symmetric* and *Hermitean*, respectively; see Sec. 3.3.

6) To have conventional definitions of algebras of pseudo-unitary and pseudo-orthogonal groups, $\hat{I}$ should commute with the matrix of *pseudo-Euclidean* metrics. This means that $\hat{I}$ should be of a *diagonal* form since symmetric form of $\hat{I}$ is not sufficient here; see Sec. 3.3.

All these requirements taken together put strong limitation on the form of $\hat{I}$, confining us with the choice made in Sec. 2.4. Namely, the family of possible units consists of *diagonal $n \times n$ matrices with positive real elements*,

$$\hat{I} = \text{diag}(q_1, q_2, \dots, q_n), \quad q_i > 0, \quad (3.79)$$

that is in confirmation with the result by Santilli.

3.5 Infinite-dimensional case

Most of the definitions and properties of $M(n, \mathbb{C}, \hat{\times})$ studied in previous sections can be readily extended to infinite dimensional case, $n = \infty$. Here, the unit is, evidently,

$$\hat{I} = \text{diag}(q_1, q_2, \dots, q_i, \dots), \quad \hat{I} \in M(\infty, \mathbb{C}, \hat{\times}), \quad (3.80)$$

and the $\hat{\times}$ -product is as usual; see Eq. (2.3). Further, in the continuous limit we have the unit

$$\hat{I}_{p'p} = \hat{I}(p')\delta(p' - p), \quad (3.81)$$

and the product,

$$(M \hat{\times} N)_{p'p} = \int dp'' dp''' M_{p'p''} \hat{T}_{p''p'''} N_{p'''p}, \quad (3.82)$$

where

$$\hat{T}_{p'p} = \hat{I}^{-1}(p')\delta(p' - p). \quad (3.83)$$

One of the applications of $M(\infty, \mathbb{C}, \hat{\times})$ and its continuous limit which would be of interest to investigate, is quantum mechanics. It is well known that quantum mechanical Hermitean operators in any representation can be given in the form of infinite-dimensional matrices.

Considering action of Lie groups on classical spaces, in Sec. 3.2, we have seen that the coordinates x^i are given with *different* weights $\sqrt{q_i}$ by the general form of unit $\hat{I}$, in contrast to equal weights, $q_i = 1$, ascribed to the coordinates by standard unit I .

Such a property is quite natural in quantum mechanics when one deals with quantum mechanical ensemble of pure states realized with different probabilities, i.e. the pure states are given with *different* weights, and thus form a *mixed state*. This formalism concerns von Neumann's density matrix and canonical ensembles. Let us consider the standard quantum mechanical definition of the density matrix,

$$\rho_{mn} = \sum P_k \bar{c}_{mk} c_{kn}, \quad (3.84)$$

where P_k are the weights, $P_k > 0$, $\sum P_k = 1$, and c_{kn} are amplitudes, and compare it with the $\hat{\times}$ -product (2.3). We see that the density matrix ρ is obtained by $\hat{\times}$ -product,

$$\rho = \bar{c} \hat{T} c, \quad (3.85)$$

of the amplitude matrices, where $\hat{T} = \text{diag}(P_1, P_2, \dots, P_k, \dots)$. In quantum mechanics, *diagonal* elements of the density matrix,

$$w_n = \rho_{nn} = \sum P_k |c_{kn}|^2, \quad (3.86)$$

is probability density to find observable in state $|n\rangle$, in the mixed ensemble. For example, in coordinate representation,

$$\rho(x, x', t) = \sum_k P_k \psi_k^*(x') \psi_k(x), \quad (3.87)$$

and

$$w(x, t) = \rho(x, x, t) = \sum_k P_k |\psi_k(x)|^2 \quad (3.88)$$

is probability density of the coordinate x , in the mixed state ensemble.

From the above sketch we see that pure states can be associated to the standard form of unit while mixed states can be associated to general form of unit (3.80), with the identification $P_k = 1/q_k$ and normalization condition $\text{Trace } \hat{T} = 1$. Then, ψ_k are seen as components of vector in space $\mathbb{R}_q^\infty$, and (3.88) is scalar product in $\mathbb{R}_q^\infty$ given by the dual metrics $g = \hat{T}$; see Eq. (3.32).

Also, there are some other examples of using weight functions in functional analysis. For example, well known space $L_{2,\rho}[a, b]$ of complex functions is a Hilbert space if one defines scalar product as [11]

$$\langle f, g \rangle = \int_a^b dx \rho(x) f(x) \bar{g}(x), \quad (3.89)$$

where the weight function $\rho(x)$ is real and positive, in the region $[a, b]$. Suppose that polynomials $p_n(x)$ constitute orthogonal system, i.e.

$$\delta_{mn} = \int_a^b dx \rho(x) p_m(x) p_n(x). \quad (3.90)$$

Then, up to constant factors, for $\rho(x) = 1$, $a = -1$, $b = 1$ we obtain Legendre polynomials, for $\rho(x) = \exp\{-x^2\}$, $a = -\infty$, $b = +\infty$ we obtain Chebyshev-Hermite polynomials, and for $\rho(x) = \exp\{-x\}$, $a = 0$, $b = +\infty$ we obtain Chebyshev-Laguerre polynomials.

The above examples are given just to stress that some elements of infinite-dimensional (discrete or continuous) case of $M(n, \mathbb{C}, \hat{\times})$, where unit is not necessarily of standard form, are well established in quantum mechanics of mixed states and in functional analysis. In both the cases, their discrete finite-dimensional limit leads to consideration of $M(n, \mathbb{C}, \hat{\times})$ equipped by unit of a general form (3.79).

In finite-dimensional case, there is the following example where inhomogeneous dilation (3.33) is explicitly used. For system of N point-particles with different masses, in three-dimensional Euclidean space we have the following Lagrangian:

$$L = \frac{1}{2} \sum_{i=1}^{3N} m_i \dot{x}_i^2 - U(x_i), \quad (3.91)$$

where $\vec{x}_k = (x_k, x_{k+1}, x_{k+2})$ are coordinates of the particles, and $m_k = m_{k+1} = m_{k+2}$, $k = 1, \dots, 3N - 3$. Introducing $x'_i = \sqrt{m_i} x_i$, one can rewrite the above Lagrangian as

$$L = \frac{1}{2} \sum_{i=1}^{3N} \dot{x}'_i{}^2 - U(x'_i). \quad (3.92)$$

The same transformation of coordinates can be used in the case of Schrödinger equation for system of N particles.

4 Examples of matrix Lie-Santilli algebras and groups

4.1 Algebra $so(3, \mathbb{R}, \hat{\times})$ and group $SO(3, \mathbb{R}, \hat{\times})$

A number of realizations of Lie-Santilli algebras and groups have been introduced and developed so far by Santilli [3, 8, 9, 4]; see also review by Santilli and Animalu [5]. In this Section we make a close consideration and review of some examples useful in physics.

We start our consideration of examples of matrix Lie-Santilli algebras and groups with the general form of unit by consideration of $so(3, \mathbb{R}, \hat{\times})$.

Let us briefly recall the usual $so(3, \mathbb{R})$ algebra. Basis elements of this algebra are 3×3 skew-symmetric matrices,

$$X_1 = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & 1 \\ 0 & -1 & 0 \end{pmatrix}, \quad X_2 = \begin{pmatrix} 0 & 0 & -1 \\ 0 & 0 & 0 \\ 1 & 0 & 0 \end{pmatrix}, \quad X_3 = \begin{pmatrix} 0 & 1 & 0 \\ -1 & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}. \quad (4.1)$$

These matrices satisfy commutation relations

$$[X_1, X_2] = X_3, \quad [X_3, X_1] = X_2, \quad [X_2, X_3] = X_1.$$

One can construct *non-skew-symmetric* matrices

$$\hat{X}_i = X_i \hat{I} \quad \text{or} \quad \hat{X}_i = \hat{I} X_i, \quad (4.2)$$

with the unit matrix

$$\hat{I} = \text{diag}(a_1, a_2, a_3) \quad (4.3)$$

identically satisfying the $\hat{\times}$ -commutation relations

$$[\hat{X}_1, \hat{X}_2]_{\hat{\times}} = \hat{X}_3, \quad [\hat{X}_3, \hat{X}_1]_{\hat{\times}} = \hat{X}_2, \quad [\hat{X}_2, \hat{X}_3]_{\hat{\times}} = \hat{X}_1. \quad (4.4)$$

Indeed, e.g., for the first $\hat{\times}$ -commutator in the above equation we have

$$\hat{X}_1 \hat{I} \hat{X}_2 - \hat{X}_2 \hat{I} \hat{X}_1 = X_1 \hat{I} X_2 - X_2 \hat{I} X_1 = (X_1 X_2 - X_2 X_1) \hat{I} = X_3 \hat{I}. \quad (4.5)$$

However, these $\hat{X}_i$ matrices can not be used to construct the group $SO(3, \mathbb{R}, \hat{\times})$, except for the trivial case $a_1 = a_2 = a_3$, which makes $\hat{X}_i$'s skew-symmetric. Indeed, only skew-symmetric matrices correspond to the orthogonal group; see Eq. (3.69).

To construct appropriate $\hat{X}_i$'s with general values of the parameters a_i , we proceed as follows. First, we calculate for ordinary X_i 's given by (4.1) the commutators $X_i \hat{I} X_k - X_k \hat{I} X_i$. They are

$$X_1 \hat{I} X_2 - X_2 \hat{I} X_1 = a_3 X_3, \quad X_3 \hat{I} X_1 - X_1 \hat{I} X_3 = a_2 X_2, \quad X_2 \hat{I} X_3 - X_3 \hat{I} X_2 = a_1 X_1.$$

Then, we use the duality property (2.21) of the $\hat{\times}$ -commutator, and find that the matrices $\hat{X}_i$ of the form

$$\hat{X}_i = \hat{I} X_i \hat{I}, \quad (4.6)$$

are skew-symmetric by construction, namely,

$$\begin{aligned} \hat{X}_1 &= \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & a_2 a_3 \\ 0 & -a_2 a_3 & 0 \end{pmatrix}, \quad \hat{X}_2 = \begin{pmatrix} 0 & 0 & -a_1 a_3 \\ 0 & 0 & 0 \\ a_1 a_3 & 0 & 0 \end{pmatrix}, \\ \hat{X}_3 &= \begin{pmatrix} 0 & a_1 a_2 & 0 \\ -a_1 a_2 & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix}, \end{aligned} \quad (4.7)$$

and satisfy the commutation relations

$$[\hat{X}_1, \hat{X}_2]_{\hat{\times}} = a_3 \hat{X}_3, \quad [\hat{X}_3, \hat{X}_1]_{\hat{\times}} = a_2 \hat{X}_2, \quad [\hat{X}_2, \hat{X}_3]_{\hat{\times}} = a_1 \hat{X}_1. \quad (4.8)$$

Corresponding matrix exponents, namely, $\hat{O}_i(t) = e^{t\hat{X}_i}$, obtained by the use of Eq. (3.56) are of the following form:

$$\hat{O}_1 = \begin{pmatrix} a_1 & 0 & 0 \\ 0 & a_2 \cos \sqrt{a_2 a_3} t & \sqrt{a_2 a_3} \sin \sqrt{a_2 a_3} t \\ 0 & -\sqrt{a_2 a_3} \sin \sqrt{a_2 a_3} t & a_3 \cos \sqrt{a_2 a_3} t \end{pmatrix}, \quad (4.9)$$

$$\hat{O}_2 = \begin{pmatrix} a_1 \cos \sqrt{a_1 a_3} t & 0 & -\sqrt{a_1 a_3} \sin \sqrt{a_1 a_3} t \\ 0 & a_2 & 0 \\ \sqrt{a_1 a_3} \sin \sqrt{a_1 a_3} t & 0 & a_3 \cos \sqrt{a_1 a_3} t \end{pmatrix}, \quad (4.10)$$

$$\hat{O}_3 = \begin{pmatrix} a_1 \cos \sqrt{a_1 a_2} t & \sqrt{a_1 a_2} \sin \sqrt{a_1 a_2} t & 0 \\ -\sqrt{a_1 a_2} \sin \sqrt{a_1 a_2} t & a_2 \cos \sqrt{a_1 a_2} t & 0 \\ 0 & 0 & a_3 \end{pmatrix}. \quad (4.11)$$

Simple but tedious algebra shows that these matrices have determinants equal to $\text{Det } \hat{I} = a_1 a_2 a_3$, and satisfy orthogonality condition (3.7). Therefore we conclude that linear combinations of these basis elements constitute the group $\text{SO}(3, \mathbb{R}, \hat{\times})$.

Action of the $\hat{O}_i$'s on vector $r = (x_1, x_2, x_3)$ is of the form $\hat{O}_i \hat{T} r$, so we present explicit forms of the matrices $\tilde{O}_i = \hat{O}_i \hat{T}$, which are of practical use, below:

$$\tilde{O}_1 = \begin{pmatrix} 1 & 0 & 0 \\ 0 & \cos \sqrt{a_2 a_3} t & \sqrt{a_2/a_3} \sin \sqrt{a_2 a_3} t \\ 0 & -\sqrt{a_3/a_2} \sin \sqrt{a_2 a_3} t & \cos \sqrt{a_2 a_3} t \end{pmatrix}, \quad (4.12)$$

$$\tilde{O}_2 = \begin{pmatrix} \cos \sqrt{a_1 a_3} t & 0 & -\sqrt{a_1/a_3} \sin \sqrt{a_1 a_3} t \\ 0 & 1 & 0 \\ \sqrt{a_3/a_1} \sin \sqrt{a_1 a_3} t & 0 & \cos \sqrt{a_1 a_3} t \end{pmatrix}, \quad (4.13)$$

$$\tilde{O}_3 = \begin{pmatrix} \cos \sqrt{a_1 a_2} t & \sqrt{a_1/a_2} \sin \sqrt{a_1 a_2} t & 0 \\ -\sqrt{a_2/a_1} \sin \sqrt{a_1 a_2} t & \cos \sqrt{a_1 a_2} t & 0 \\ 0 & 0 & 1 \end{pmatrix}. \quad (4.14)$$

4.2 The groups $\text{SO}(2, \mathbb{R}, \hat{\times})$ and $\text{U}(1, \mathbb{C}, \hat{\times})$

According to the results of Sec. 4.1, elements of the group $\text{SO}(2, \mathbb{R}, \hat{\times})$ are of the form

$$\hat{O} = \begin{pmatrix} a_1 \cos \sqrt{a_1 a_2} t & \sqrt{a_1 a_2} \sin \sqrt{a_1 a_2} t \\ -\sqrt{a_1 a_2} \sin \sqrt{a_1 a_2} t & a_2 \cos \sqrt{a_1 a_2} t \end{pmatrix}. \quad (4.15)$$

This matrix is counterpart of the matrix of usual rotation of Euclidean plane.

We are interested in the group $\text{U}(1, \mathbb{C}, \hat{\times})$, from which, by *making it real*, the group $\text{O}(2, \mathbb{R}, \hat{\times})$ can be obtained. We recall that in the usual case elements

of $U(1)$ are complex numbers of unit module, e^{it} . The representation (4.15) can be reproduced by the following *realization map*.

First, we note that for the 2×2 unit matrix $\hat{I} = \text{diag}(a_1, a_2)$, we have $\text{Det } \hat{I} = a_1 a_2$, and $\text{Det } \hat{O} = a_1 a_2 = \text{Det } \hat{I}$, as it should be for special orthogonal matrices. We introduce the complex number

$$\xi = \text{Det } \hat{I} \exp\{i\sqrt{\text{Det } \hat{I}} t\}, \quad (4.16)$$

and the matrix

$$\hat{J} = \sqrt{\text{Det } \hat{I}} \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}. \quad (4.17)$$

Observe that $|\xi| = \text{Det } \hat{I}$, and the matrix $\hat{J}$ has the properties

$$\text{Det } \hat{J} = \text{Det } \hat{I}, \quad \hat{J}^2 \equiv \hat{J}\hat{T}\hat{J} = -\hat{I}, \quad (4.18)$$

and does not commute with $\hat{I}$,

$$\hat{I}\hat{J} - \hat{J}\hat{I} = \sqrt{\text{Det } \hat{I}}(a_1 - a_2) \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}. \quad (4.19)$$

However, it $\hat{\times}$ -commute with $\hat{I}$, namely,

$$\hat{I}\hat{\times}\hat{J} - \hat{J}\hat{\times}\hat{I} = 0. \quad (4.20)$$

Also, note that, particularly, $\hat{O}(0) = \hat{I}$ and $\hat{O}(\pi/2) = \hat{J}$.

Then, it can be verified that the realization map is given by

$$\hat{O} = r(\xi) \equiv \text{Re } \xi (\text{Det } \hat{T})\hat{I} + \text{Im } \xi (\text{Det } \hat{T})\hat{J}. \quad (4.21)$$

Indeed, $\text{Re } \xi = \text{Det } \hat{I} \cos \sqrt{a_1 a_2} t$, $\text{Im } \xi = \text{Det } \hat{I} \sin \sqrt{a_1 a_2} t$, and multiplying these by $(\text{Det } \hat{T})\hat{I}$ and $(\text{Det } \hat{T})\hat{J}$, respectively, we reproduce, after summing up, the representation (4.15). Note that determinant of the realization map matrix is

$$\text{Det } r(\xi) = |\xi|, \quad (4.22)$$

and matrix r $\hat{\times}$ -commutes with $\hat{J}$ due to Eq. (4.20) that means that $\hat{J}$ is indeed an operator of complex character.

Thus, elements of $U(1, \mathbb{C}, \hat{\times})$ are of the form (4.16), with product of the complex numbers given trivially by

$$\xi_3 = \xi_1 \hat{\times} \xi_2 \equiv \xi_1 (\text{Det } \hat{T}) \xi_2, \quad \xi_{1,2,3} \in U(1, \mathbb{C}, \hat{\times}), \quad (4.23)$$

where $\text{Det } \hat{T} = 1/(a_1 a_2)$ is a real number, and a_1 and a_2 are fixed positive real numbers.

Note that the matrix product in $SO(2, \mathbb{R}, \hat{\times})$ is $\hat{O}_1 \hat{T} \hat{O}_2$ while in $U(1, \mathbb{C}, \hat{\times})$ the product is due to the above rule, where $\text{Det } \hat{T}$ is used instead of $\hat{T}$. Due to the realization map (4.21), the groups $SO(2, \mathbb{R}, \hat{\times})$ and $U(1, \mathbb{C}, \hat{\times})$ are isomorphic to each other.

An important remark here is that in the realization map (4.21) we used the fact that $-\hat{J}\hat{T}\hat{J} = \hat{I}$ due to Eq. (4.18). By this, we achieved isomorphism between $\text{SO}(2, \mathbb{R}, \hat{\times})$ and $\text{U}(1, \mathbb{C}, \hat{\times})$. Indeed, one can see from the form (4.16) of ξ that the group $\text{U}(1, \mathbb{C}, \hat{\times})$ is characterized by *one* independent parameter $a = \text{Det } \hat{I} = a_1 a_2$, while $\text{SO}(2, \mathbb{R}, \hat{\times})$ is characterized by *two* independent parameters a_1 and a_2 . So, if we were used the matrix $-\hat{J}^2 = \hat{J}\hat{J} = (\text{Det } \hat{I})I$ as a unit matrix we would obtain $\text{SO}(2, \mathbb{R}, \hat{\times})$ characterized by the only parameter a instead of the two parameters a_1 and a_2 . Namely, the unit would be of trivialized form $a_1 a_2 \text{diag}(1, 1)$ and elements of group $\text{SO}(2, \mathbb{R}, \hat{\times})$ would be of the form

$$\hat{O} = \sqrt{a_1 a_2} \begin{pmatrix} \cos \sqrt{a_1 a_2} t & \sin \sqrt{a_1 a_2} t \\ -\sin \sqrt{a_1 a_2} t & \cos \sqrt{a_1 a_2} t \end{pmatrix}. \quad (4.24)$$

Thus, the lesson is that we should use $-\hat{J}\hat{T}\hat{J}$ rather than $-\hat{J}\hat{J}$ to define the unit matrix in realization map. Of course, we have some features stemming from the real dimensionality two. See Sec. 4.5 for general consideration of the realization map, $n > 2$.

For convenience and to have consistence with the general definition (3.9), we take $\xi \in \text{U}(1, \mathbb{C}, \hat{\times})$ in the form (4.16). Note that we can replace $\text{Det } \hat{I}$ by any real parameter but we are using $\text{Det } \hat{I}$ to keep explicit correspondence with $\text{SO}(2, \mathbb{R}, \hat{\times})$.

In fact, it does not matter which non-zero value the *module* of ξ has because it can be absorbed by appropriate definition of the product (4.23) and associated realization map (4.21). For example, we can put $\xi = \sqrt{\text{Det } \hat{I}} \exp\{i\sqrt{\text{Det } \hat{I}} t\}$ provided that $\text{Det } \hat{T}$ in Eqs. (4.23) and (4.21) is replaced by $\sqrt{\text{Det } \hat{T}}$, obtaining the same result. Moreover, we can put simply $\xi = \exp\{i\sqrt{\text{Det } \hat{I}} t\}$, i.e. $|\xi| = 1$, and, accordingly, drop $\text{Det } \hat{T}$ in Eqs. (4.23) and (4.21).

This demonstrates the fact that group $\text{U}(1, \mathbb{C}, \hat{\times})$ is isomorphic to usual $\text{U}(1)$, up to the factor $\sqrt{\text{Det } \hat{I}}$ in the argument of complex number. This factor is of importance since $\sqrt{\text{Det } \hat{I}}$ is fixed, and

$$\exp\{i\sqrt{\text{Det } \hat{I}} t_1\} \exp\{i\sqrt{\text{Det } \hat{I}} t_2\} = \exp\{i\sqrt{\text{Det } \hat{I}} (t_1 + t_2)\} \quad (4.25)$$

is again element of $\text{U}(1, \mathbb{C}, \hat{\times})$ for any t_1 and t_2 , while, for example,

$$\exp\{i\sqrt{\text{Det } \hat{I}} t_1\} \exp\{it_2\} \quad (4.26)$$

is not element of the group for any t_1 and t_2 .

In other words, $\text{U}(1, \mathbb{C}, \hat{\times})$ consists of complex numbers ξ with modules $|\xi| = \text{Det } \hat{I}$ and arguments $\text{Arg } \xi$ dividable to $\sqrt{\text{Det } \hat{I}}$, i.e. arguments modulo number $\sqrt{\text{Det } \hat{I}}$.

4.3 Action of the group $\text{U}(1, \mathbb{C}, \hat{\times})$

Let us consider the action of group $\text{U}(1, \mathbb{C}, \hat{\times})$.

While $U(1)$ linearly transforms $\mathbb{C}$ equipped by standard metrics

$$|z|^2 = x^2 + y^2,$$

the group $U(1, \mathbb{C}, \hat{\times})$ must conserve, by definition (3.2), metrics

$$|z|^{\hat{2}} = a_1 x^2 + a_2 y^2,$$

which is *not conformally equivalent* to the standard metrics $|z|^2$. In fact, we see that group $U(1) \subset \mathbb{C}$ and produces motion of $\mathbb{C} = \mathbb{R}^2$ while group $U(1, \mathbb{C}, \hat{\times}) \subset \mathbb{C}$ and produces, with the action defined by (3.19), motion of $\mathbb{C}_q = \mathbb{R}_q^2$, where $\mathbb{R}_q^2$ is Euclidean (flat) space equipped by the metrics $a_1 x^2 + a_2 y^2$.

The fact that we can rescale module of $\xi \in U(1, \mathbb{C}, \hat{\times})$ to 1 without loss of generality corresponds to the conformal equivalence of the metrics $a_1 x^2 + a_2 y^2$ and $b(a_1 x^2 + a_2 y^2)$, where b is a real constant.

We are interested to find out transformation of $\mathbb{C}_q$ corresponding to rotation (4.15) of the space $\mathbb{R}_q^2$.

Let us consider action of $SO(2, \mathbb{R}, \hat{\times})$ on $\mathbb{R}_q^2$. Action of (4.15) on vector $r = (x, y)$ reads $\hat{O}\hat{\times}r$, namely,

$$\hat{O}\hat{T}r = \begin{pmatrix} \cos \sqrt{a_1 a_2} t & \sqrt{a_1/a_2} \sin \sqrt{a_1 a_2} t \\ -\sqrt{a_2/a_1} \sin \sqrt{a_1 a_2} t & \cos \sqrt{a_1 a_2} t \end{pmatrix} \begin{pmatrix} x \\ y \end{pmatrix}. \quad (4.27)$$

The matrix $\tilde{O} = \hat{O}\hat{T}$ has the following particular values:

$$\tilde{O}(0) = I, \quad \tilde{O}(\pi/2) = \hat{J}\hat{T} = (\text{Det } \hat{T})\hat{I}\hat{J}. \quad (4.28)$$

The transformation (4.27) is a linear one,

$$\begin{pmatrix} a & b \\ -c & a \end{pmatrix} \begin{pmatrix} x \\ y \end{pmatrix}, \quad (4.29)$$

where a , b , and c are real parameters. However, linear transformation of complex space, $z \mapsto \lambda z$, where $\lambda = (a + ib)$, $z = (x + iy) \in \mathbb{C}$, results in

$$\begin{pmatrix} a & b \\ -b & a \end{pmatrix} \begin{pmatrix} x \\ y \end{pmatrix}. \quad (4.30)$$

Obviously, the use of modified transformation $z \mapsto \lambda \hat{\times} z = \lambda (\text{Det } \hat{T})z$ does not yield transformation of the form (4.29) because this causes just additional dilation by the real factor $\text{Det } \hat{T}$.

Thus, we are led to consider transformation of the general form

$$z \mapsto F(z, \bar{z}).$$

Let us consider standard $\mathbb{R}^2$ and make the inhomogeneous dilation (3.33) of its coordinates,

$$x' = x/\sqrt{a_1}, \quad y' = y/\sqrt{a_2}. \quad (4.31)$$

Then, $r^2 = x^2 + y^2$ becomes

$$r^2 = a_1 x'^2 + a_2 y'^2 = r' \hat{T} r' = r'^2. \quad (4.32)$$

Since $a_{1,2} > 0$ the transformation (4.31) is invertible and well defined. Associated Jacobi matrix is the same as the transformation matrix of (4.31), namely,

$$\begin{pmatrix} 1/\sqrt{a_1} & 0 \\ 0 & 1/\sqrt{a_2} \end{pmatrix} = +\sqrt{\hat{T}} \quad (4.33)$$

and Jacobian is $\text{Det } \hat{T}$. This transformation obviously provides the map $\mathbb{R}^2 \mapsto \mathbb{R}_q^2$ due to Eq. (4.32); $r = (x, y) \in \mathbb{R}^2$, $r' = (x', y') \in \mathbb{R}_q^2$.

In terms of complex coordinates, using $x = (z + \bar{z})/2$, $y = (z - \bar{z})/2i$ we obtain from the transformation (4.31)

$$z' = \frac{z + \bar{z}}{2\sqrt{a_1}} + \frac{z - \bar{z}}{2\sqrt{a_2}}, \quad \bar{z}' = \frac{z + \bar{z}}{2\sqrt{a_1}} - \frac{z - \bar{z}}{2\sqrt{a_2}}, \quad (4.34)$$

or

$$z' = \left(\frac{1}{2\sqrt{a_1}} + \frac{1}{2\sqrt{a_2}}\right)z + \left(\frac{1}{2\sqrt{a_1}} - \frac{1}{2\sqrt{a_2}}\right)\bar{z} \equiv f(z, \bar{z}), \quad (4.35)$$

$$\bar{z}' = \left(\frac{1}{2\sqrt{a_1}} + \frac{1}{2\sqrt{a_2}}\right)z - \left(\frac{1}{2\sqrt{a_1}} - \frac{1}{2\sqrt{a_2}}\right)\bar{z} = f(z, -\bar{z}). \quad (4.36)$$

Function f in the transformation (4.35) depends on $\bar{z}$, and thus it is *not* complex-analytic function ($\partial f / \partial \bar{z} \neq 0$), which thus makes *non complex-analytic transformation* of complex plane $\mathbb{C}$. We write in this case $\mathbb{C} \mapsto \mathbb{C}_q$, to avoid confusion with the usual convention that *transformation* of complex plane, $\mathbb{C} \mapsto \mathbb{C}$, means *complex-analytic* transformation. Accordingly, we write $z' \in \mathbb{C}_q$.

Function $f(z, \bar{z})$ is a sum of holomorphic and antiholomorphic functions, $f(z, \bar{z}) = f_1(z) + f_2(\bar{z})$, each of which is a linear function of its argument.

Below, we make various linear transformations of complex plane $\mathbb{C}$, namely, $\mathbb{C} \mapsto \mathbb{C}$, and analyze what kind of transformations they induce in complex plane $\mathbb{C}_q$, namely, $\mathbb{C}_q \mapsto \mathbb{C}_q$.

Let us make linear complex-analytic transformation of $\mathbb{C}$, namely, $z \mapsto \lambda z$, and see its image in $\mathbb{C}_q$. We observe from (4.35) that $f_1(z)$ becomes $\lambda f_1(z)$ while $f_2(\bar{z})$ remains intact, and thus we have *no* linear transformation of $\mathbb{C}_q$ which is of the form $z' \mapsto \lambda z'$. By construction, this describes the action of group $U(1)$, for $|\lambda| = 1$. So, the image of $U(1)$ -action on $\mathbb{C}$ is some *non-linear* transformation in $\mathbb{C}_q$.

By making linear *non* complex-analytic transformation of $\mathbb{C}$, namely, $z \mapsto \lambda(z + \bar{z})$, we readily obtain that the image of this transformation in $\mathbb{C}_q$ is *linear complex-analytic* transformation, $z' \mapsto \lambda z'$, of $\mathbb{C}_q$. However, this still does not describe the action of group $U(1, \mathbb{C}, \hat{\times})$.

By making linear *non* complex-analytic transformation of $\mathbb{C}$, namely, $z \mapsto \lambda(z + \bar{z}) + \mu(z - \bar{z})$, we obtain the image of this transformation in $\mathbb{C}_q$ which is *linear non complex-analytic* transformation of $\mathbb{C}_q$ which do correspond to the action of group $U(1, \mathbb{C}, \hat{\times})$.

Indeed, by choosing complex numbers in the form

$$\lambda = \frac{1}{2}(\sqrt{a_1}a + i\sqrt{a_2}b), \quad \mu = \frac{i}{2}(\sqrt{a_1}b - i\sqrt{a_2}a), \quad (4.37)$$

where a and b are arbitrary real numbers, and making non complex-analytic transformation of $\mathbb{C}_q$ as above,

$$z' \mapsto \lambda(z' + \bar{z}') + \mu(z' - \bar{z}'), \quad (4.38)$$

we obtain directly the following associated transformation of $\mathbb{R}_q^2$:

$$\begin{pmatrix} x' \\ y' \end{pmatrix} \mapsto \begin{pmatrix} a & \sqrt{a_2/a_1}b \\ -\sqrt{a_1/a_2}b & a \end{pmatrix} \begin{pmatrix} x' \\ y' \end{pmatrix}, \quad (4.39)$$

which is of the required form (4.29).

4.4 Group $\text{SO}(1, 1, \mathbb{R}, \hat{\times})$

In the usual setting of group $\text{SO}(1, 1, \mathbb{R})$, we have the generator of the form

$$X = \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}, \quad (4.40)$$

and the usual matrix exponent gives us elements of the proper group in the form

$$O = \begin{pmatrix} \cosh \phi & \sinh \phi \\ \sinh \phi & \cosh \phi \end{pmatrix}. \quad (4.41)$$

This matrix is pseudo-orthogonal, $OGO^t = G$, where $G = \text{diag}(1, -1)$, and $\text{Det } O = 1$, so that the scalar product $xGx = (x^0)^2 - (x^1)^2$ is conserved; x^0 and x^1 are local coordinates of two-dimensional pseudo-Euclidean (Minkowski) space M^2 . Also, there is a smooth path from O to $I = \text{diag}(1, 1)$.

In the group $\text{SO}(1, 1, \mathbb{R}, \hat{\times})$, the unit is $\hat{I} = \text{diag}(a_0, a_1)$; $\hat{T} = (\hat{I})^{-1}$. The generator can be chosen here as $\hat{X} = \hat{I}X\hat{I}$ (see Sec. 4.1),

$$\hat{X} = a_0a_1 \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}, \quad (4.42)$$

which supplies us, by the help of the matrix exponent $\hat{e}^{\phi X} = e^{\phi X \hat{T} \hat{I}}$, with elements of the group having the form (cf. Eq. (4.15))

$$\hat{O} = \begin{pmatrix} a_0 \cosh \sqrt{a_0 a_1} \phi & \sqrt{a_0 a_1} \sinh \sqrt{a_0 a_1} \phi \\ \sqrt{a_0 a_1} \sinh \sqrt{a_0 a_1} \phi & a_1 \cosh \sqrt{a_0 a_1} \phi \end{pmatrix}. \quad (4.43)$$

This matrix conserves the scalar product

$$x\hat{G}x = a_0(x^0)^2 - a_1(x^1)^2, \quad (4.44)$$

where metrics is

$$\hat{G} = \hat{I}G = \begin{pmatrix} a_0 & 0 \\ 0 & -a_1 \end{pmatrix}, \quad (4.45)$$

in the sense that

$$\hat{O} \hat{\times} \hat{G} \hat{\times} \hat{O}^t \equiv \hat{O} \hat{T} \hat{G} \hat{T} \hat{O}^t = \hat{G}. \quad (4.46)$$

Also, $\hat{O}$ can be continuously connected to the identity transformation $\hat{I}$ by $\phi \rightarrow 0$. It is instructive to check the above pseudo-orthogonality (4.46), where $\hat{T} = \text{diag}(1/a_0, 1/a_1)$ is used for the product, since $\hat{O}(\psi) \hat{\times} \hat{I} \hat{\times} \hat{O}^t(\psi) \neq \hat{I}$, in contrast to the expectation that $\hat{O}(\psi) \hat{\times} \hat{I} \hat{\times} \hat{O}^t(\psi) = \hat{I}$. In fact, this is equal to $\hat{O}$ with *double* angle, $\hat{O}(\psi) \hat{\times} \hat{I} \hat{\times} \hat{O}^t(\psi) = \hat{O}(2\psi)$. Moreover, similar (simple but tedious) calculations show that in addition to Eq. (4.46), we have

$$\hat{O} \hat{W} \hat{G} \hat{W} \hat{O}^t = \hat{G}, \quad (4.47)$$

where we have denoted

$$\hat{W} = \hat{G}^{-1} = \begin{pmatrix} 1/a_0 & 0 \\ 0 & -1/a_1 \end{pmatrix}. \quad (4.48)$$

Our comment here is that we can use the pair $(\hat{G}, \hat{W})$ instead of $(\hat{I}, \hat{T})$.

Action of the matrix $\hat{O}$ on $r = (x^0, x^1)$ is $\hat{O} \hat{T} r$, namely,

$$x'^0 = x^0 \cosh \psi + x^1 \sqrt{\frac{a_0}{a_1}} \sinh \psi, \quad x'^1 = \sqrt{\frac{a_1}{a_0}} x^0 \sinh \psi + x^1 \cosh \psi, \quad (4.49)$$

where we have denoted $\psi = \sqrt{a_0 a_1} \phi$, for brevity. At $x^1 = 0$, we have from Eq. (4.49)

$$\frac{x'^1}{x'^0} = \sqrt{\frac{a_1}{a_0}} \tanh \psi. \quad (4.50)$$

In the context of *special relativity in two dimensions*, $x^0 = ct$, the l.h.s. of Eq. (4.50) is the relative speed v/c of the frame of reference (ct, x^1) with respect to the frame of reference (ct', x'^1) . Thus,

$$\tanh \psi = \sqrt{\frac{a_0}{a_1}} \frac{v}{c}. \quad (4.51)$$

Inserting this into Eq. (4.49), we obtain by the use of trigonometric relations $\sinh \psi = \tanh \psi / \sqrt{1 - (\tanh \psi)^2}$ and $\cosh \psi = 1 / \sqrt{1 - (\tanh \psi)^2}$,

$$t' = t \frac{1}{\sqrt{1 - \hat{\beta}^2}} + x^1 \frac{a_0}{a_1} \frac{v/c^2}{\sqrt{1 - \hat{\beta}^2}}, \quad x'^1 = t \frac{v}{\sqrt{1 - \hat{\beta}^2}} + x^1 \frac{1}{\sqrt{1 - \hat{\beta}^2}}, \quad (4.52)$$

where we have denoted

$$\hat{\beta} = \sqrt{\frac{a_0}{a_1}} \frac{v}{c}. \quad (4.53)$$

Since only the ratio a_0/a_1 is present in Eq. (4.52), we denote

$$a = \sqrt{\frac{a_1}{a_0}} \quad (4.54)$$

and rewrite it in a more compact form

$$t' = (t + \frac{v}{a^2 c^2} x^1) \hat{\gamma}, \quad x'^1 = (x^1 + vt) \hat{\gamma}, \quad (4.55)$$

where we denote the $\hat{\gamma}$ -factor

$$\hat{\gamma} = \frac{1}{\sqrt{1 - \hat{\beta}^2}} \quad (4.56)$$

and $\hat{\beta} = v/(ac)$.

The following remarks are in order.

(a) We took without proof that the trigonometric relations used above in the case of metrics $\text{diag}(a_1, -a_2)$ are the same as they are in the case of standard pseudo-Euclidean metrics $\text{diag}(1, -1)$, and note only that the spaces are flat in both the cases.

(b) Despite the fact that $\hat{I}$ depends on two parameters a_0 and a_1 , and the generator $\hat{X}$ depends on the product $a_1 a_2$, only their *ratio* (4.54) has appeared in the transformations (4.55).

(c) We see from Eq. (4.55) that the only distinction from the conventional Lorentz transformations is that the constant c is replaced by ac . This can be understood as follows. Making inhomogeneous dilation (rescaling) of the coordinates $x^0 \mapsto x^0/\sqrt{a_0}$ and $x^1 \mapsto x^1/\sqrt{a_1}$ to obtain metrics G from $\hat{G}$, we simply change the slope of isotropic line, $x^1 = ct$ to $x^1 = act$.

(d) We expect that such properties extend to the action of the higher-dimensional pseudo-orthogonal group $\text{SO}(3, 1, \mathbb{R}, \hat{\times})$ on the corresponding four-dimensional Minkowski space-time, with three different coefficients appearing at c in three main space axeses Ox^1 , Ox^2 , and Ox^3 of a chosen coordinate system (space anisotropic behavior).

4.5 Realization map

The realization map constructed in Sec. 4.2 can be extended to higher dimensions in the following way.

First, we note that $\text{GL}(n, \mathbb{C}, \hat{\times})$ and $\text{GL}(m, \mathbb{R}, \hat{\times})$, $m = 2n$, in general have the units parameterized by n and $2n$ parameters, respectively. In the two-dimensional case it appeared fortunately that $-\hat{J}^2 = \hat{I}$ exactly for special choice of the parameters. However, this is not the case in higher dimensions, $n > 2$, and some reparameterization is needed to achieve match between the parameters.

Indeed, the $2n \times 2n$ matrix $\hat{J}$ has the form

$$\hat{J} = \begin{pmatrix} 0 & \hat{I} \\ -\hat{I} & 0 \end{pmatrix}. \quad (4.57)$$

where $n \times n$ matrix $\hat{I} = \text{diag}(q_1, q_2, \dots, q_n)$ is unit in $\text{GL}(n, \mathbb{C}, \hat{\times})$. Note that $\text{Det } \hat{J} = (\text{Det } \hat{I})^2$. The corresponding $2n \times 2n$ unit $\hat{I}_{2n}$ in $\text{GL}(2n, \mathbb{R}, \hat{\times})$ can be found by squaring the matrix $\hat{J}$ with the help of unit matrix $\hat{I}_{gen}$,

$$\hat{I}_{gen} = \text{diag}(a_1, a_2, \dots, a_{2n}), \quad (4.58)$$

of $\text{GL}(m, \mathbb{R}, \hat{\times})$, $m = 2n$. Namely,

$$\hat{I}_{2n} = -\hat{J} \hat{\times} \hat{J} = -\hat{J} \hat{T}_{gen} \hat{J} = \begin{pmatrix} K_1 & 0 \\ 0 & K_2 \end{pmatrix}, \quad (4.59)$$

where $\hat{T}_{gen} = \hat{I}_{gen}^{-1}$ and

$$K_1 = \text{diag}(q_1^2/a_{n+1}, q_2^2/a_{n+2}, \dots, q_n^2/a_{2n}), \quad (4.60)$$

$$K_2 = \text{diag}(q_1^2/a_1, q_2^2/a_2, \dots, q_n^2/a_n). \quad (4.61)$$

Then, the unit matrix $\hat{I}_{2n}$ depends on $2n$ independent parameters, with extra n independent parameters coming from $\hat{I}_{gen}$. The realization map for matrix $M = A + iB$, $M \in \text{GL}(n, \mathbb{C}, \hat{\times})$, is given by

$$r(M) = \hat{I}_{2n} A + \hat{J} B, \quad r(M) \in \text{GL}(2n, \mathbb{C}, \hat{\times}). \quad (4.62)$$

However, $\hat{I}_{2n}$ is not equal to $\hat{I}_{gen}$ even if we identify $q_i = a_i$, $i = 1, \dots, n$. Explicitly, in this case we have

$$K_1 = \text{diag}(a_1^2/a_{n+1}, a_2^2/a_{n+2}, \dots, a_n^2/a_{2n})$$

and

$$K_2 = \text{diag}(a_1, a_2, \dots, a_n).$$

Only after the reparametrization

$$a_i^2/a_{n+i} \mapsto a_{n+i}, \quad i = 1, \dots, n, \quad (4.63)$$

we obtain $K_1 = \text{diag}(a_{n+1}, a_{n+2}, \dots, a_{2n})$, and thus achieve the identification $\hat{I}_{2n} = \hat{I}_{gen}$.

Let us consider the case $n = 4$ for an illustrative purpose. Let $\text{GL}(2, \mathbb{C}, \hat{\times})$ and $\text{GL}(4, \mathbb{R}, \hat{\times})$ have the units

$$\hat{I} = \begin{pmatrix} q_1 & 0 \\ 0 & q_2 \end{pmatrix}, \quad \hat{I}_{gen} = \begin{pmatrix} a_1 & 0 & 0 & 0 \\ 0 & a_2 & 0 & 0 \\ 0 & 0 & a_3 & 0 \\ 0 & 0 & 0 & a_4 \end{pmatrix}, \quad (4.64)$$

respectively, and $\hat{T}_{gen} = \hat{I}_{gen}^{-1}$. Then, $\hat{J}$ is

$$\hat{J} = \begin{pmatrix} 0 & 0 & q_1 & 0 \\ 0 & 0 & 0 & q_2 \\ -q_1 & 0 & 0 & 0 \\ 0 & -q_2 & 0 & 0 \end{pmatrix}, \quad (4.65)$$

and unit $\hat{I}_4 = \hat{J}\hat{T}_{gen}\hat{J}$ has the form

$$\hat{I}_4 = \begin{pmatrix} q_1^2/a_3 & 0 & 0 & 0 \\ 0 & q_2^2/a_4 & 0 & 0 \\ 0 & 0 & q_1^2/a_1 & 0 \\ 0 & 0 & 0 & q_2^2/a_2 \end{pmatrix}. \quad (4.66)$$

Putting $q_1 = a_1$ and $q_2 = a_2$, and reparameterizing $a_1^2/a_3 \mapsto a_3$ and $a_2^2/a_4 \mapsto a_4$ we reproduce $\hat{I}_4 = \hat{I}_{gen}$.

4.6 Matrix algebra $M(2, \mathbb{C}, \hat{\times})$

Let us consider $M(2, \mathbb{C})$ matrix algebra consisting of all 2×2 matrices over the field of complex numbers $\mathbb{C}$.

Additive basis of the matrix algebra $M(2, \mathbb{C})$ consists of unit 2×2 matrix I and matrices σ_i ,

$$\sigma_1 = \begin{pmatrix} i & 0 \\ 0 & -i \end{pmatrix}, \quad \sigma_2 = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}, \quad \sigma_3 = \begin{pmatrix} 0 & i \\ i & 0 \end{pmatrix}. \quad (4.67)$$

This means that the algebra with the basis

$$\{I, \sigma_1, \sigma_2, \sigma_3\}, \quad (4.68)$$

and relations

$$\sigma_i \sigma_j + \sigma_j \sigma_i = -2I \delta_{ij}, \quad (4.69)$$

over $\mathbb{C}$, i.e. the universal enveloping algebra of $su(2, \mathbb{C})$, is isomorphic to $M(2, \mathbb{C})$. Note that the above σ -matrices are traceless and skew-Hermitian in $M(2, \mathbb{C})$, and are related to usual Pauli matrices by factor i , with labels 1 and 3 interchanged.

Following Santilli, we define new unit matrix

$$\hat{I} = \begin{pmatrix} q_1 & 0 \\ 0 & q_2 \end{pmatrix}, \quad q_{1,2} > 0, \quad \hat{I} \in M(2, \mathbb{C}), \quad (4.70)$$

and the associated $\hat{\times}$ -product between the matrices

$$M \hat{\times} N = M \hat{T} N, \quad M, N \in M(2, \mathbb{C}), \quad (4.71)$$

where

$$\hat{T} = \hat{I}^{-1} = \begin{pmatrix} 1/q_1 & 0 \\ 0 & 1/q_2 \end{pmatrix}, \quad \hat{T} \in M(2, \mathbb{C}). \quad (4.72)$$

Explicitly,

$$M \hat{\times} N = M \hat{T} N = \begin{pmatrix} m_{11} & m_{12} \\ m_{21} & m_{22} \end{pmatrix} \begin{pmatrix} q_1^{-1} & 0 \\ 0 & q_2^{-1} \end{pmatrix} \begin{pmatrix} n_{11} & n_{12} \\ n_{21} & n_{22} \end{pmatrix} \quad (4.73)$$

$$= \begin{pmatrix} \frac{m_{11}n_{11}}{q_1} + \frac{m_{12}n_{21}}{q_2} & \frac{m_{11}n_{12}}{q_1} + \frac{m_{12}n_{22}}{q_2} \\ \frac{m_{21}n_{11}}{q_1} + \frac{m_{22}n_{21}}{q_2} & \frac{m_{21}n_{12}}{q_1} + \frac{m_{22}n_{22}}{q_2} \end{pmatrix}. \quad (4.74)$$

We would like to construct an additive basis

$$\{\hat{I}, i\hat{\sigma}_1, i\hat{\sigma}_2, i\hat{\sigma}_3\}, \quad (4.75)$$

in terms of which elements of the algebra $M(2, \mathbb{C}, \hat{\times})$ are presented as *linear combinations*. Namely,

$$M = x_0\hat{I} + x_1i\hat{\sigma}_1 + x_2i\hat{\sigma}_2 + x_3i\hat{\sigma}_3, \quad M \in M(2, \mathbb{C}, \hat{\times}), \quad (4.76)$$

where x_i are parameters. Note that we are not using $\hat{\times}$ -product to multiply matrices by parameters (numbers) in Eq. (4.76) because parameters are not elements of the matrix algebra, and we are not considering action of matrices on a vector.

The criterium to determine $\hat{\sigma}$ -matrices is that they, together with the unit $\hat{I}$, must form additive basis in algebra $M(2, \mathbb{C}, \hat{\times})$.

One of the possible ways to construct such a basis is that $\hat{\sigma}$ -matrices must satisfy the following anticommutation relations, instead of the standard (4.69),

$$\hat{\sigma}_i\hat{T}\hat{\sigma}_j + \hat{\sigma}_j\hat{T}\hat{\sigma}_i = 2\hat{I}\delta_{ij}, \quad (4.77)$$

or, defining $\hat{\times}$ -anticommutator,

$$\{M, N\}_{\hat{\times}} = M\hat{\times}N + N\hat{\times}M = M\hat{T}N + N\hat{T}M, \quad (4.78)$$

we rewrite the above as

$$\{\hat{\sigma}_i, \hat{\sigma}_j\}_{\hat{\times}} = 2\hat{I}\delta_{ij}. \quad (4.79)$$

We have two formal algebraic solutions for these equations,

$$\hat{\sigma}_i = \sigma_i\hat{I}, \quad (4.80)$$

i.e.

$$\hat{\sigma}_1 = \begin{pmatrix} iq_1 & 0 \\ 0 & -iq_2 \end{pmatrix}, \quad \hat{\sigma}_2 = \begin{pmatrix} 0 & q_2 \\ -q_1 & 0 \end{pmatrix}, \quad \hat{\sigma}_3 = \begin{pmatrix} 0 & iq_2 \\ iq_1 & 0 \end{pmatrix}, \quad (4.81)$$

and

$$\hat{\sigma}_i = \hat{I}\sigma_i, \quad (4.82)$$

i.e.

$$\hat{\sigma}_1 = \begin{pmatrix} iq_1 & 0 \\ 0 & -iq_2 \end{pmatrix}, \quad \hat{\sigma}_2 = \begin{pmatrix} 0 & q_1 \\ -q_2 & 0 \end{pmatrix}, \quad \hat{\sigma}_3 = \begin{pmatrix} 0 & iq_1 \\ iq_2 & 0 \end{pmatrix}. \quad (4.83)$$

Indeed, we have identically for (4.82)

$$\hat{\sigma}_i\hat{T}\hat{\sigma}_j + \hat{\sigma}_j\hat{T}\hat{\sigma}_i = \hat{I}\sigma_i\hat{T}\hat{I}\sigma_j + \hat{I}\sigma_j\hat{T}\hat{I}\sigma_i = \hat{I}(\sigma_i\sigma_j + \sigma_j\sigma_i) = -\hat{I}2\delta_{ij}, \quad (4.84)$$

and similarly for (4.80).

Note that the relations (4.84) hold for any invertible 2×2 matrix $\hat{I}$, not only for those having the above mentioned diagonal form (4.70). This means that the algebraic solutions (4.80) and (4.82) are formal.

We stress that the spaces $M(2, \mathbb{C})$ and $M(2, \mathbb{C}, \hat{\times})$ are isomorphic to each other. The difference is that they have different bases, $B = \{I, \sigma_1, \sigma_2, \sigma_3\}$ and $\hat{B} = \{\hat{I}, \hat{\sigma}_1, \hat{\sigma}_2, \hat{\sigma}_3\}$, respectively, and different definitions of matrix product. In view of the solutions (4.80) and (4.82), these bases are related to each other simply by

$$\hat{B} = B\hat{I}, \quad (4.85)$$

and

$$\hat{B} = \hat{I}B, \quad (4.86)$$

respectively. From this point of view, if one change the unit I by some transformation matrix, the same matrix should be used to change remaining elements of the basis. This justifies partially the choice of algebraic solutions in the forms (4.80) and (4.82).

An important note is that in general matrices $\hat{\sigma}_i$ are not skew-Hermitian and not traceless. However, due to Eq. (3.58), we must have $\text{Trace } \hat{\sigma}_i \hat{I} = 0$ for matrices (4.80) to meet the condition $|\text{Det } U| = \text{Det } \hat{I}$ for the associated Lie group. This is indeed trivially the case. The problem of the lack of skew-Hermiticity concerns algebra $su(2, \mathbb{C}, \hat{\times})$, and is considered in the following Section.

4.7 Algebra $su(2, \mathbb{C}, \hat{\times})$ and group $SU(2, \mathbb{C}, \hat{\times})$

Let us find norm of the vector

$$X = x_1 i \hat{\sigma}_1 + x_2 i \hat{\sigma}_2 + x_3 i \hat{\sigma}_3 \quad (4.87)$$

in space $su(2, \mathbb{C}, \hat{\times})$. Using Eq. (4.81) and noting from Eq. (4.77) that $\hat{\sigma}_1^2 = -\hat{I}$, $\hat{\sigma}_2^2 = -\hat{I}$, and $\hat{\sigma}_3^2 = -\hat{I}$, we have (the Killing metrics)

$$|X|^2 = \text{Det } X = q_1 q_2 (x_1^2 + x_2^2 + x_3^2) = (\text{Det } \hat{I})(x_1^2 + x_2^2 + x_3^2). \quad (4.88)$$

Transformation

$$X \mapsto u \hat{\times} X \hat{\times} u^{-1}, \quad u \in su(2, \mathbb{C}, \hat{\times}), \quad (4.89)$$

is orthogonal in the sense of the scalar product (4.88), namely,

$$\text{Det } X = \text{Det } (u \hat{\times} X \hat{\times} u^{-1}). \quad (4.90)$$

Thus, any matrix $Z \in su(2, \mathbb{C}, \hat{\times})$ makes linear transformation $\text{ad } Z = [Z, X]_{\hat{\times}}$ of three-dimensional space $su(2, \mathbb{C}, \hat{\times})$.

Metric tensor of space $su(2, \mathbb{C}, \hat{\times})$ due to Eq. (4.88) is, evidently,

$$\hat{\delta}_{ij} = \delta_{ij} \text{Det } \hat{I}, \quad (4.91)$$

which is conformally equivalent to the usual Euclidean metrics δ_{ij} of three-dimensional Euclidean space $\mathbb{R}^3$.

Formally, $\hat{\sigma}$ -matrices form representation of algebra $su(2, \mathbb{C}, \hat{\times})$. To prove this, we must verify the following commutation relations:

$$[\hat{\sigma}_1, \hat{\sigma}_2]_{\hat{\times}} = 2\hat{\sigma}_3, \quad [\hat{\sigma}_3, \hat{\sigma}_1]_{\hat{\times}} = 2\hat{\sigma}_2, \quad [\hat{\sigma}_2, \hat{\sigma}_3]_{\hat{\times}} = 2\hat{\sigma}_1. \quad (4.92)$$

One can easily verify by using (4.80) or (4.82) that these relations trivially hold.

Note, however, that in general the matrices $\hat{\sigma}_i$ are *not skew-Hermitean*. This give us no possibility to construct associated *unitary* group with their help.

They become skew-symmetric in the particular case, $q_1 = q_2 = q$, that leads, however, to reduction of $\hat{I}$ to scalar matrix

$$\hat{I} = \text{diag}(q, q) = qI, \quad (4.93)$$

and therefore trivializes the attempt.

Also, application of the duality method developed for $so(3, \mathbb{R}, \hat{\times})$ in Sec. 4.1 to the case of $su(2, \mathbb{C}, \hat{\times})$ does not seem to provide us with an appropriate algebraic solution. The obstacle is made by σ_1 matrix, which has a *diagonal* form whereas none of X_i 's has a diagonal form. Explicit calculations show that

$$\begin{aligned} \sigma_1 \hat{I} \sigma_2 - \sigma_2 \hat{I} \sigma_1 &= (q_1 + q_2) \sigma_3, & \sigma_1 \hat{I} \sigma_3 - \sigma_3 \hat{I} \sigma_1 &= -(q_1 + q_2) \sigma_2, \\ \sigma_2 \hat{I} \sigma_3 - \sigma_3 \hat{I} \sigma_2 &= q_1 q_2 \hat{T} \sigma_1, \end{aligned} \quad (4.94)$$

implying that the matrices

$$\hat{\sigma}_i = \hat{I} \sigma_i \hat{I} \quad (4.95)$$

having explicit *skew-Hermitean* form

$$\hat{\sigma}_1 = \begin{pmatrix} iq_1^2 & 0 \\ 0 & -iq_2^2 \end{pmatrix}, \quad \hat{\sigma}_2 = \begin{pmatrix} 0 & q_1 q_2 \\ -q_1 q_2 & 0 \end{pmatrix}, \quad \hat{\sigma}_3 = \begin{pmatrix} 0 & iq_1 q_2 \\ iq_1 q_2 & 0 \end{pmatrix}, \quad (4.96)$$

satisfy

$$[\hat{\sigma}_1, \hat{\sigma}_2]_{\hat{\times}} = (q_1 + q_2) \hat{\sigma}_3, \quad [\hat{\sigma}_1, \hat{\sigma}_3]_{\hat{\times}} = -(q_1 + q_2) \hat{\sigma}_2, \quad [\hat{\sigma}_2, \hat{\sigma}_3]_{\hat{\times}} = q_1 q_2 \hat{T} \hat{\sigma}_1. \quad (4.97)$$

Here, the last equation includes the matrix $\hat{T}$ so that these commutation relations are seemed to be not Lie-algebraic. Also, direct calculations show that $\hat{\times}$ -anticommutators between these $\hat{\sigma}$ -matrices are of the form

$$\{\hat{\sigma}_1, \hat{\sigma}_1\}_{\hat{\times}} = -2\hat{I}^3, \quad \{\hat{\sigma}_2, \hat{\sigma}_2\}_{\hat{\times}} = -2q_1 q_2 \hat{I}, \quad \{\hat{\sigma}_3, \hat{\sigma}_3\}_{\hat{\times}} = -2q_1 q_2 \hat{I}, \quad (4.98)$$

$$\{\hat{\sigma}_1, \hat{\sigma}_2\}_{\hat{\times}} = iq_1 q_2 (q_1 - q_2) \sigma_2, \quad \{\hat{\sigma}_1, \hat{\sigma}_3\}_{\hat{\times}} = iq_1 q_2 (q_1 - q_2) \sigma_3, \quad \{\hat{\sigma}_2, \hat{\sigma}_3\}_{\hat{\times}} = 0. \quad (4.99)$$

We note that some non-zero values appear in Eq. (4.99).

Thus, the problem to construct general solution for $\hat{\sigma}$ -matrices which obeys appropriate $\hat{\times}$ -anticommutation and/or $\hat{\times}$ -commutation relations, and are *skew-Hermitean*, is opened. Note that we have explicit construction for $so(3, \mathbb{R}, \hat{\times})$ with general values of the parameters a_i , $i = 1, 2, 3$; see Sec. 4.1. And this is a candidate to the algebra isomorphic to $su(2, \mathbb{C}, \hat{\times})$, with unrestricted parameters q_i , $i = 1, 2$. However, we should to note that number of the parameters a_i and q_i is different.

Elements of the group $SU(2, \mathbb{C}, \hat{\times})$ can be represented by using the matrix exponent (see Sec. 3.3),

$$M = \hat{e}^{\frac{1}{2} t^i \hat{\sigma}_i}, \quad M \in SU(2, \mathbb{C}, \hat{\times}), \quad (4.100)$$

where t^i are real parameters.

Direct calculations show that for the representations (4.80) and (4.82) we obtain matrix exponents, which indeed exhibit the property $\text{Det } M = \text{Det } \hat{I}$, but they are *not unitary*. Evidently the latter is a consequence of the fact that these representations are not skew-Hermitean matrices. In the case of unit of the form of scalar matrix (4.93), the group $\text{SU}(2, \mathbb{C}, \hat{\times})$ is simply isomorphic to the ordinary group $\text{SU}(2, \mathbb{C})$ since q can be absorbed by the parameters t^i .

In the case of the representation (4.96), we have $\text{Trace } \hat{\sigma}_1 \hat{T} = (q_1 - q_2)$, $\text{Trace } \hat{\sigma}_2 \hat{T} = 0$, $\text{Trace } \hat{\sigma}_1 \hat{T} = 0$, that means that $\text{Det } M \neq \text{Det } \hat{I}$ for $M = \hat{e}^{\frac{1}{2}t^1 \hat{\sigma}_1}$, and thus we can not construct group $\text{SU}(2, \mathbb{C}, \hat{\times})$ despite the fact that these $\hat{\sigma}$ -matrices are skew-Hermitean. Instead, we could construct $\text{U}(2, \mathbb{C}, \hat{\times})$ but only if the $\hat{\times}$ -commutation relations (4.97) are acceptable.

A Appendix A

Distributivity implies that abstract product $f(M, N)$ is a *linear* function in both the matrices,

$$f(M_1 + M_2, N) = f(M_1, N) + f(M_2, N),$$

$$f(M, N_1 + N_2) = f(M, N_1) + f(M, N_2),$$

restricting possible functions $f(M, N)$ by a polynomial in M and N . Let us define the product in the form

$$f(M, N) = \tau_1 M \tau_2 + \tau_3 N \tau_4 + \tau_5 M \tau_6 N \tau_7 + \tau_8 N \tau_9 M \tau_{10}, \quad (\text{A.1})$$

where τ_i are fixed matrices; $M, N, \tau_i \in \mathcal{M}(n, \mathbb{C})$. Axiom of left and right unit gives us two equations

$$f(\hat{I}, N) = \hat{I},$$

$$f(N, \hat{I}) = N,$$

which should be *solvable* equations for any N to provide us the algebra with unit. Namely, for (A.1) we have

$$\tau_1 \hat{I} \tau_2 + \tau_3 N \tau_4 + \tau_5 \hat{I} \tau_6 N \tau_7 + \tau_8 N \tau_9 \hat{I} \tau_{10} = N, \quad (\text{A.2})$$

$$\tau_1 N \tau_2 + \tau_3 \hat{I} \tau_4 + \tau_5 N \tau_6 \hat{I} \tau_7 + \tau_8 \hat{I} \tau_9 N \tau_{10} = N, \quad (\text{A.3})$$

from which we see that to satisfy identically the equations each term in the l.h.s. of them must be considered separately. Namely,

$$\tau_1 \hat{I} \tau_2 = 0, \quad (\text{A.4})$$

$$\tau_3 = I, \quad \tau_4 = I, \quad (\text{A.5})$$

$$\tau_5 \hat{I} \tau_6 = I, \quad \tau_7 = I, \quad (\text{A.6})$$

$$\tau_8 = I, \quad \tau_9 \hat{I} \tau_{10} = I, \quad (\text{A.7})$$

$$\tau_3 \hat{I} \tau_4 = 0, \quad (\text{A.8})$$

$$\tau_1 = I, \quad \tau_2 = I, \quad (\text{A.9})$$

$$\tau_5 = I, \quad \tau_6 \hat{I} \tau_7 = I, \quad (\text{A.10})$$

$$\tau_8 \hat{I} \tau_9 = I, \quad \tau_{10} = I. \quad (\text{A.11})$$

Therefore, since we assume $\hat{I} \neq 0$ we have $\tau_1 = 0$ or $\tau_2 = 0$ and $\tau_3 = 0$ or $\tau_4 = 0$, that rules out first two terms in (A.1). Note that the same result can be obtained by using the distributivity condition. Further, we obtain

$$\hat{I} = \tau_6^{-1}, \quad \text{or} \quad \hat{I} = \tau_9^{-1}, \quad (\text{A.12})$$

where we have assumed that τ_6 and τ_9 are invertible matrices. This means that we are left with the following two forms of product,

$$f(M, N) = M \tau_6 N \quad \text{or} \quad f(M, N) = N \tau_9 M, \quad (\text{A.13})$$

which are in essence equivalent to each other.

Putting the constant terms $\tau_1 \hat{I} \tau_2$ and $\tau_3 \hat{I} \tau_4$ to zero is an obvious requirement, while putting the remaining terms of Eqs. (A.2) and (A.3) *separately* equal to N needs some comments. To see more closely on the above made separation of the terms let us check the associativity condition, $f(f(M, N), P) = f(M, f(N, P))$,

$$\begin{aligned} f(f(M, N), P) &= \tau_5 (\tau_5 M \tau_6 N \tau_7 + \tau_8 N \tau_9 M \tau_{10}) \tau_6 P \tau_7 \\ &\quad + \tau_8 P \tau_9 (\tau_5 M \tau_6 N \tau_7 + \tau_8 N \tau_9 M \tau_{10}) \tau_{10}, \end{aligned} \quad (\text{A.14})$$

$$\begin{aligned} f(M, f(N, P)) &= \tau_5 M \tau_6 (\tau_5 N \tau_6 P \tau_7 + \tau_8 P \tau_9 N \tau_{10}) \tau_7 \\ &\quad + \tau_8 (\tau_5 N \tau_6 P \tau_7 + \tau_8 P \tau_9 N \tau_{10}) \tau_9 M \tau_{10}. \end{aligned} \quad (\text{A.15})$$

Obviously, the associativity condition is not satisfied for this general form of the product. The term $\tau_5 \tau_8 N \tau_9 M \tau_{10} \tau_6 P \tau_7$ is present in Eq. (A.14) while such a term is absent in Eq. (A.15). Therefore, to meet the associativity condition we must put one of fixed matrices in this term equal to zero. This leads to discarding either third ($\tau_5 = 0$ or $\tau_6 = 0$ or $\tau_7 = 0$) or fourth ($\tau_8 = 0$ or $\tau_9 = 0$ or $\tau_{10} = 0$) term in the definition (A.1), thus yielding its separate consideration. For $\tau_8 = 0$ or $\tau_9 = 0$ or $\tau_{10} = 0$, the associativity condition reads

$$\tau_5 \tau_5 M \tau_6 N \tau_7 \tau_6 P \tau_7 = \tau_5 M \tau_6 \tau_5 N \tau_6 P \tau_7 \tau_7, \quad (\text{A.16})$$

from which we find again $\tau_5 = \tau_7 = I$. For $\tau_5 = 0$ or $\tau_6 = 0$ or $\tau_7 = 0$, we find similarly $\tau_8 = \tau_{10} = I$. As the conclusion, we obtain the product in the form (2.3).

For completeness, let us consider higher degrees (> 2) in M or N in definition of the product. Using axiom of left and right unit, $f(\hat{I}, N) = f(N, \hat{I}) = N$, one can see that there is no possibility to have these equations identically satisfied for fixed τ_i and any N . For example, the definition

$$f(M, N) \doteq \tau_1 M \tau_2 N \tau_3 N + N \tau_4 M \tau_5 N \tau_6 + N \tau_7 N \tau_8 M \tau_9, \quad (\text{A.17})$$

where τ_i are fixed matrices, $M, N, \tau_i \in \mathcal{M}(n, \mathbb{C})$, implies

$$\tau_1 \hat{I} \tau_2 N \tau_3 N + N \tau_4 \hat{I} \tau_5 N \tau_6 + N \tau_7 N \tau_8 \hat{I} \tau_9 = N, \quad (\text{A.18})$$

$$\tau_1 N \tau_2 \hat{I} \tau_3 \hat{I} + \hat{I} \tau_4 N \tau_5 \hat{I} \tau_6 + \hat{I} \tau_7 \hat{I} \tau_8 N \tau_9 = N \quad (\text{A.19})$$

These two equations can not be identically satisfied for fixed τ_i and arbitrary N . Indeed, in the first equation, matrix N appears two times in each term of the l.h.s. so that some of τ_i must be of the form N^{-1} to satisfy this equation. However, we assume that τ_i 's are *fixed* matrices so that they can not be of the form N^{-1} , where N is an *arbitrary* matrix. The same reason rules out any higher degree in M or N . Thus, the form (2.3) is the most general form for associative and distributive product in matrix algebra with unit.

B Appendix B

Below, we present the proof of the statement that matrices I and $\hat{I}$ are not similar to each other. It is based on the construction of the associated invariant polynomials [2].

Let us consider the matrices of the form $(I\lambda - A)$, where A is 2×2 matrix and λ is a real number.

Two matrices A and B are *similar* to each other iff the matrices $X = (I\lambda - A)$ and $Y = (I\lambda - B)$ have the same *invariant polynomials*.

For the case under study, $A = I$ and $B = \hat{I}$, we have

$$X = (\lambda - 1)I = \text{diag}(\lambda - 1, \lambda - 1). \quad (\text{B.1})$$

and

$$Y = (I\lambda - \hat{I}) = \text{diag}(\lambda - q_1, \lambda - q_2). \quad (\text{B.2})$$

Let us find invariant polynomials of X . Main minors of X are

$$\begin{aligned} \text{2nd order minors:} & \quad (\lambda - 1)^2, \\ \text{1st order minors:} & \quad (\lambda - 1), \quad (\lambda - 1). \end{aligned} \quad (\text{B.3})$$

Largest common quotients of the minors are

$$d_2(\lambda) = (\lambda - 1)^2, \quad d_1(\lambda) = (\lambda - 1). \quad (\text{B.4})$$

Then, the invariant polynomials are

$$i_2(\lambda) \equiv d_2/d_1 = (\lambda - 1), \quad i_1(\lambda) \equiv d_1 = (\lambda - 1). \quad (\text{B.5})$$

Let us find invariant polynomials of Y . Main minors of Y are

$$\begin{aligned} \text{2nd order minors:} & \quad (\lambda - q_1)(\lambda - q_2), \\ \text{1st order minors:} & \quad (\lambda - q_1), \quad (\lambda - q_2). \end{aligned} \quad (\text{B.6})$$

Then, largest common quotients of the minors of are

$$d_2(\lambda) = (\lambda - q_1)(\lambda - q_2), \quad d_1(\lambda) = 1. \quad (\text{B.7})$$

Thus, invariant polynomials are

$$i_2(\lambda) = d_2/d_1 = (\lambda - q_1)(\lambda - q_2), \quad i_1(\lambda) = d_1 = 1. \quad (\text{B.8})$$

We see that the invariant polynomials of I given by (B.5) and of $\hat{I}$ given by (B.8) are different. Thus, the matrices I and $\hat{I}$ are not similar to each other in the sense of (2.60) and (2.61).

Perhaps, the above presented strict proof is not necessary to see that I and $\hat{I}$ are not similar to each other since it is almost obvious. However, we have seen that I and $\hat{I}$ are homotopically equivalent in the space of matrices, for $q_{1,2} > 0$, that could be thought of as they are related to each other by the similarity condition $I = V^{-1}\hat{I}V$ for some matrix V . We have proved explicitly that this is not the case.

One can easily prove that this property holds for general n -dimensional case, $M(n, \mathbb{C})$, by noting that $n = 2$ case forms subspace of the higher-dimensional cases.

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